合
institut POLYTECHNIQUE DE PARIS

Alexis Drouot<br>Ubiquity of conical points in topological insulators<br>Tome 8 (202I), p. 5o7-532.<br><http://jep.centre-mersenne.org/item/JEP_2021_<br>$\qquad$ 8 507_0>

© Les auteurs, 2021.
Certains droits réservés.

(c) (1)
Cet article est mis à disposition selon les termes de la licence
Licence internationale d'attribution Creative Commons BY 4.0.
https://creativecommons.org/licenses/by/4.0/
L'accès aux articles de la revue «Journal de l'École polytechnique - Mathématiques » (http://jep.centre-mersenne.org/), implique l'accord avec les conditions générales d'utilisation (http://jep.centre-mersenne.org/legal/).


MERSENNE

# UBIQUITY OF CONICAL POINTS IN TOPOLOGICAL INSULATORS 

by Alexis Drouot


#### Abstract

We show that generically, the degeneracies of a family of Hermitian matrices depending on three parameters have a conical structure. Our result applies to the study of topological phases of matter. It suggests that adiabatic deformations of two-dimensional topological insulators come generically with Dirac-like propagating currents, whose total conductivity equals the chiral number of conical points. Résumé (Omniprésence des points de Dirac dans les isolants topologiques) Nous montrons que les valeurs propres dégénérées de matrices dépendant de trois paramètres possèdent généralement une structure conique. Nous appliquons ce résultat à l'étude des phases topologiques de systèmes quantiques. Nous montrons que les déformations adiabatiques entre deux isolants topologiques distincts ont une conductivité globale égale au nombre chiral de points de Dirac.


## Contents

1. Introduction ..... 507
2. $\mathbb{M}$ is open ..... 515
3. $\mathbb{M}$ is dense ..... 517
4. Proof of Theorem 1 ..... 524
5. Chern number difference. ..... 525
Appendix. Continuous approximation. ..... 529
References ..... 530

## 1. Introduction

Let $\mathcal{E}$ be the space of $N \times N$ Hermitian matrices; $\mathcal{E}^{*} \subset \mathcal{E}$ consisting of matrices with simple eigenvalues; and $\mathbb{T}^{2}$ be a two-dimensional torus. Given $H_{0}$ and $H_{1}$ in $C^{\infty}\left(\mathbb{T}^{2}, \varepsilon^{*}\right)$, is there a path from $H_{0}$ to $H_{1}$, that remains in $C^{\infty}\left(\mathbb{T}^{2}, \varepsilon^{*}\right) ?$

[^0]In general, the response is no: there is a topological obstruction, related to the eigenbundles of $H_{0}$ and $H_{1}$. When this obstruction is present, any path from $H_{0}$ to $H_{1}$ acquires degenerate eigenvalues. In this paper, we explore the shape of these crossings. We show that generically, they exhibit a conical structure.

This result has a counterpart in topological phases of matter. When two topologically distinct insulators are adiabatically connected, it suggests that generically:

- Finitely many channels supporting chiral currents appear;
- Up to large times, these currents follow a Dirac equation and are concentrated (in phase-space) along conical eigenvalue crossings;
- The chiral number of currents equals the Chern number difference.

This establishes a quantitative link between (a) asymmetric currents; (b) eigenvalue crossings; and (c) the bulk-edge correspondence.
1.1. Genericity of conical points. - We first state our result in a form that applies to topological phases of matter. We postpone the general statement to Section 1.4.

Let $\mathbb{T}^{2}=\mathbb{R}^{2} /(2 \pi \mathbb{Z})^{2}$ and $H_{0}, H_{1}$ be two elements of $C^{\infty}\left(\mathbb{T}^{2}, \mathcal{E}\right)$, with eigenvalues $\lambda_{1}\left(H_{j}(\xi)\right) \leqslant \cdots \leqslant \lambda_{N}\left(H_{j}(\xi)\right)$, repeated according to multiplicity. We assume that for some $n \in[1, N-1]$ and all $\xi \in \mathbb{T}^{2}$,

$$
\begin{equation*}
\lambda_{n}\left(H_{0}(\xi)\right)<\lambda_{n+1}\left(H_{0}(\xi)\right), \quad \lambda_{n}\left(H_{1}(\xi)\right)<\lambda_{n+1}\left(H_{1}(\xi)\right) . \tag{1.1}
\end{equation*}
$$

Let $\mathcal{L}$ be the set of smooth homotopies from $H_{0}$ to $H_{1}$ :

$$
\mathcal{L} \stackrel{\text { def }}{=}\left\{H \in C^{\infty}\left([0,1] \times \mathbb{T}^{2}, \mathcal{E}\right): H(0, \cdot)=H_{0}, H(1, \cdot)=H_{1}\right\} .
$$

Definition 1. - If $H \in \mathcal{L}$, we say that $\lambda_{n}(H)$ and $\lambda_{n+1}(H)$ cross (or degenerate) at $\zeta_{0}=\left(s_{0}, \xi_{0}\right) \in[0,1] \times \mathbb{T}^{2}$ if $\lambda_{n}\left(H\left(\zeta_{0}\right)\right)=\lambda_{n+1}\left(H\left(\zeta_{0}\right)\right)$.

We say that $\lambda_{n}(H)$ and $\lambda_{n+1}(H)$ cross conically if $\lambda_{n}\left(H\left(\zeta_{0}\right)\right)$ has multiplicity precisely two; and if there exist $a_{0} \in \mathbb{R}^{3}$ and $S_{0} \in M_{3}(\mathbb{R})$ invertible such that

$$
\left\{\begin{array}{c}
\lambda_{n}\left(H\left(\zeta_{0}+\varepsilon\right)\right)=\lambda_{n}\left(H\left(\zeta_{0}\right)\right)+\left\langle a_{0}, \varepsilon\right\rangle-\left\|S_{0} \varepsilon\right\|+o(\varepsilon)  \tag{1.2}\\
\lambda_{n+1}\left(H\left(\zeta_{0}+\varepsilon\right)\right)=\lambda_{n}\left(H\left(\zeta_{0}\right)\right)+\left\langle a_{0}, \varepsilon\right\rangle+\left\|S_{0} \varepsilon\right\|+o(\varepsilon)
\end{array}, \quad \varepsilon \in \mathbb{R}^{3} \quad\right. \text { small. }
$$

Conical degeneracies correspond to tilted cones in the graphs of eigenvalues - see Figure 1. In particular, conical crossings of $\lambda_{n}(H)$ and $\lambda_{n+1}(H)$ are isolated. At first, we could think that they are rare among degeneracies: a non-empty intersection of two surfaces is in general a curve (rather than a point). Nonetheless:

Theorem 1. - If $H_{0}$ and $H_{1}$ are elements of $C^{\infty}\left(\mathbb{T}^{2}, \varepsilon\right)$ satisfying (1.1), then

$$
\mathbb{L}=\left\{H \in \mathcal{L}: \text { all crossings of } \lambda_{n}(H) \text { and } \lambda_{n+1}(H) \text { are conical }\right\}
$$

is a dense open subset of $\mathcal{L}$.
The natural topology on $\mathcal{L}$ is that induced by $C^{\infty}\left([0,1] \times \mathbb{T}^{2}, \mathcal{E}\right)$, see Section 1.7. Results at lower regularity are also possible; our techniques typically require $C^{2}$.

As a simple consequence of Theorem 1, for generic $H \in \mathcal{L}, \lambda_{n}(H)$ and $\lambda_{n+1}(H)$ cross at only finitely many points: conical crossings are isolated. Under a topological


Figure 1. (a) Eigenvalue surfaces of $H\left(s_{0}, \cdot\right)$ near a conical point $\left(s_{0}, \xi_{0}\right)$ of $H$. They intersect at the vertex of a (non-isotropic) cone. (b) Eigenvalue surfaces of $H(s, \cdot)$ for $s \neq s_{0}$ near $s_{0}$. They no longer touch.
condition on $H_{0}$ and $H_{1}$, crossings must nonetheless arise. Indeed, (1.1) allows us to define a rank- $n$ vector bundle $\mathscr{V}_{0}$ over $\mathbb{T}^{2}$ : the fibers are

$$
\mathscr{V}_{0}(\xi)=\bigoplus_{j=1}^{n} \operatorname{ker}\left(H_{0}(\xi)-\lambda_{j}\left(H_{0}(\xi)\right)\right), \quad \xi \in \mathbb{T}^{2}
$$

We can also define $\mathscr{V}_{1}$, associated to $H_{1}$ : only (1.1) is necessary to construct such vector bundles. Hence, if there is a homotopy between $H_{0}$ and $H_{1}$ that maintains (1.1), then there are smooth vector bundles $\mathscr{V}_{s} \rightarrow \mathbb{T}^{2}, s \in[0,1]$, interpolating between $\mathscr{V}_{0}$ and $\mathscr{V}_{1}$. In particular, $\mathscr{V}_{0}$ and $\mathscr{V}_{1}$ would be topologically equivalent.

This restriction can be measured via the Chern number - the vector bundle analog of the Euler characteristic. This number can take any integer value, even in the context of eigenbundles - see the appendix in [Dro19c] - and characterizes the topology when the basis is a two-torus - see e.g. [Pan07, Mon17]. Thus, $\mathscr{V}_{0}$ and $\mathscr{V}_{1}$ are topologically equivalent if and only if $c_{1}\left(\mathscr{V}_{0}\right)=c_{1}\left(\mathscr{V}_{1}\right)$. In particular, if $H_{0}, H_{1} \in C^{\infty}(\mathcal{M}, \mathcal{E})$ satisfy (1.1) and $c_{1}\left(\mathscr{V}_{0}\right) \neq c_{1}\left(\mathscr{V}_{1}\right)$, then any homotopy between $H_{0}$ and $H_{1}$ admits degeneracies. These, according to Theorem 1, are generically all conical - see Figure 2.
1.2. Connection with topological phases of matter. - We review tight-binding, translation-invariant models of insulators at an energy $\lambda_{0} \in \mathbb{R}$. These systems are represented by selfadjoint Hamiltonians $\mathcal{H}_{0}: \ell^{2}\left(\mathbb{Z}^{2}, \mathbb{C}^{N}\right) \rightarrow \ell^{2}\left(\mathbb{Z}^{2}, \mathbb{C}^{N}\right)$ with:

$$
\begin{equation*}
\left[\mathcal{H}_{0}, T_{j}\right]=0, \quad\left(T_{j} \psi\right)_{m}=\psi_{m+e_{j}} ; \quad \text { and } \quad \lambda_{0} \notin \sigma\left(\mathcal{H}_{0}\right) \tag{1.3}
\end{equation*}
$$

In (1.3), $\sigma\left(\mathcal{H}_{0}\right)$ denotes the $\ell^{2}\left(\mathbb{Z}^{2}, \mathbb{C}^{N}\right)$-spectrum of $\mathcal{H}_{0}$. Physically, $\lambda_{0} \notin \sigma\left(\mathcal{H}_{0}\right)$ means that there is no plane-wave propagation at energy $\lambda_{0}$.

Thanks to (1.3) and $\left[T_{1}, T_{2}\right]=0$, we can diagonalize $\mathcal{H}_{0}, T_{1}$ and $T_{2}$ simultaneously. The eigenvalues of $T_{j}$ are $e^{i \xi_{j}}, \xi_{j} \in \mathbb{T}^{1}=\mathbb{R} /(2 \pi \mathbb{Z})$. Joint eigenspaces of $T_{1}$ and $T_{2}$


Figure 2. For each $n \in[1, N-1], C^{\infty}\left(\mathbb{T}^{2}, \mathcal{E}\right)$ splits in components distinguished by Chern numbers. If $H_{0}$ and $H_{1}$ lie in different components, a path joining $H_{0}$ to $H_{1}$ (blue) acquire crossings. Non-conicaltype degeneracies (red) are rare in $C^{\infty}\left(\mathbb{T}^{2}, \mathcal{E}\right)$.
canonically identify with $\mathbb{C}^{N}$ :

$$
\bigcap_{j=1}^{2} \operatorname{ker}\left(T_{j}-e^{i \xi_{j}}\right)=\left\{\left(e^{i \xi m} \psi_{0}\right)_{m \in \mathbb{Z}^{2}}: \psi_{0} \in \mathbb{C}^{N}\right\}, \quad \xi=\left(\xi_{1}, \xi_{2}\right) \in \mathbb{T}^{2}
$$

Thus, the analysis of $\mathcal{H}_{0}$ reduces to that of its Bloch transform: the $\mathbb{T}^{2}$-parametrized family of $N \times N$ Hermitian matrices

$$
\begin{equation*}
H_{0}(\xi)=e^{-i \xi m} \cdot \mathcal{H}_{0} \cdot e^{i \xi m}, \quad \xi \in \mathbb{T}^{2} \tag{1.4}
\end{equation*}
$$

The insulating condition $\lambda_{0} \notin \sigma\left(\mathcal{H}_{0}\right)$ and the spectral decomposition of $\mathcal{H}_{0}$ into $\left\{H_{0}(\xi)\right\}_{\xi \in \mathbb{T}^{2}}$ imply that $\lambda_{0}$ is never in $\sigma\left(H_{0}(\xi)\right)$. Thus, $H_{0}$ satisfies (1.1).

A standard question in topological phases of matter is whether two materials can be deformed to each other while maintaining their electronic properties. If $\mathcal{H}_{1}$ is another insulator at energy $\lambda_{0}$, with associated vector bundle $\mathscr{V}_{1}$ of rank $n$, then $H_{1}$ also satisfies (1.1). As explained in Section 1.1, if $c_{1}\left(\mathscr{V}_{0}\right) \neq c_{1}\left(\mathscr{V}_{1}\right)$, then there are no path $\left\{\mathcal{H}_{s}\right\}_{s \in[0,1]}$ connecting $\mathcal{H}_{0}$ and $\mathcal{H}_{1}$ while maintaining (1.3). Physically, two topologically distinct insulators cannot be deformed to one another without passing by a conductor.

Theorem 1 explains quantitatively this failure. Generically, conical crossings arise as one transitions from $\mathcal{H}_{0}$ to $\mathcal{H}_{1}$. The quantity $c_{1}\left(\mathscr{V}_{1}\right)-c_{1}\left(\mathscr{V}_{0}\right)$ is fundamental in the analysis of interface effects between topological insulators; see e.g. [RH08, Bal19a, Dro19c]. Below, we express it as the number of conical crossings, counted according to chirality.

Assume that $H \in \mathbb{L}$ and $\lambda_{n}(H)$ and $\lambda_{n+1}(H)$ degenerate conically at $\zeta_{0}$; and define $(J f)_{j}=\left\langle f, f_{j}\right\rangle$, where $\left(f_{1}, f_{2}\right)$ is an orthogonal basis of $\operatorname{ker}\left(\lambda_{n}\left(\zeta_{0}\right)-H\left(\zeta_{0}\right)\right)$. We write a Taylor expansion of the $2 \times 2$ matrix $J H\left(\zeta_{0}+\varepsilon\right) J^{*}$ near $\varepsilon=0$ :

$$
\begin{equation*}
J H\left(\zeta_{0}+\varepsilon\right) J^{*}=J H\left(\zeta_{0}\right) J^{*}+\sum_{j=1}^{3}\left(A_{0} \varepsilon\right)_{j} \cdot \sigma_{j}+O\left(\varepsilon^{2}\right) \tag{1.5}
\end{equation*}
$$

where $\sigma_{1}, \sigma_{2}, \sigma_{3}$ are the standard Pauli matrices and $A_{0} \in M_{3}(\mathbb{R})$. Using the conical structure, $A_{0}$ is invertible - see (5.7) below. The quantity $\operatorname{sgn}\left(\operatorname{det}\left(A_{0}\right)\right)$ is called the chirality of the conical point. For such degeneracies, it coincides with the topological Weyl charge defined in [MP14].

Theorem 2. - Let $H \in \mathbb{L}$, such that $\lambda_{n}(H)$ and $\lambda_{n+1}(H)$ degenerate conically precisely at $\zeta_{1}, \ldots, \zeta_{K}$. If $\operatorname{sgn}\left(\operatorname{det}\left(A_{1}\right)\right), \ldots, \operatorname{sgn}\left(\operatorname{det}\left(A_{K}\right)\right)$ are the associated chiralities, then

$$
\begin{equation*}
c_{1}\left(\mathscr{V}_{1}\right)-c_{1}\left(\mathscr{V}_{0}\right)=\sum_{k=1}^{K} \operatorname{sgn}\left(\operatorname{det}\left(A_{k}\right)\right) \tag{1.6}
\end{equation*}
$$

Theorem 1 guarantees that $\mathbb{L} \neq \varnothing$ - in fact, that $\mathbb{L}$ is a residual set.

### 1.3. Relation with adiabatic transport and bulk-edge correspondence

In this section, we explain the physical consequences of Theorems 1 and 2 on transport in adiabatic deformations of topological insulators.

Let $\mathcal{H}_{0}$ and $\mathcal{H}_{1}$ be two Hamiltonians satisfying (1.3). Let $\left\{\mathcal{H}_{s}\right\}_{s \in[0,1]}$ be a homotopy between $\mathcal{H}_{0}$ and $\mathcal{H}_{1}$; extend $\mathcal{H}_{s}$ by $\mathcal{H}_{0}$ for $s \leqslant 0$ and by $\mathcal{H}_{1}$ for $s \geqslant 1$. For $\delta>0$, we define Hamiltonians $\mathbb{Q}^{\delta}$ and $\mathcal{H}^{\delta}$ on $\ell^{2}\left(\mathbb{Z}^{2}, \mathbb{C}^{N}\right)$ by

$$
\begin{equation*}
\left(\mathbb{Q}^{\delta} \psi\right)_{m}=\left(\mathcal{H}_{\delta m_{2}} \psi\right)_{m}, \quad m=\left(m_{1}, m_{2}\right) \in \mathbb{Z}^{2} ; \quad \mathcal{H}^{\delta}=\frac{\mathbb{Q}^{\delta}+\left(\mathbb{Q}^{\delta}\right)^{*}}{2} \tag{1.7}
\end{equation*}
$$

Both $Q^{\delta}$ and $\mathcal{H}^{\delta}$ model a (spatial) deformation from $\mathcal{H}_{0}$ to $\mathcal{H}_{1}$ transversely to $\mathbb{R} e_{1}$, occurring at speed $\delta$. In addition, $\mathcal{H}^{\delta}$ is selfadjoint; and for small $\delta, \mathcal{H}^{\delta}-Q^{\delta}=O(\delta)$ as operators on $\ell^{2}\left(\mathbb{Z}^{2}, \mathbb{C}^{N}\right)$.

We are interested in the adiabatic scaling: $\delta \rightarrow 0$. This regime has an important place in the mathematical physics literature; see e.g. [Sim83, Ber84, PST03, FT16]. It corresponds to changing $\mathcal{H}_{0}$ to $\mathcal{H}_{1}$ globally (i.e., on a scale $\delta^{-1} \gg 1$ ) while preserving translation-invariance locally (i.e., on a scale $\delta^{-1 / 2}$ - note $1 \ll \delta^{-1 / 2} \ll \delta^{-1}$ ).

Generically, $\lambda_{n}\left(H_{\delta n_{2}}\right)$ and $\lambda_{n+1}\left(H_{\delta n_{2}}\right)$ do not degenerate for most values of $\delta n_{2}$. For such values, we can define the local Chern number of $\mathcal{H}^{\delta}$ at $\left(n_{1}, \delta n_{2}\right)$ : it is that of $H_{\delta n_{2}}$. The local Chern number is discontinuous at degeneracies, see Figure 3.

In adiabatic domain-wall deformations of honeycomb structures, edge states arise and are concentrated near Dirac points (isotropic conical points) [FLTW16, LTWZ19, Dro19b, DW20]. At leading order, they propagate according to an emerging Dirac operator, in the direction prescribed by chirality.

The analysis of [FLTW16, LTWZ19, Dro19b, DW20] is local in nature and would extend beyond Dirac points. Theorem 1 shows that degeneracies are generically conical. Hence, the Dirac-type propagation of edge states is universal in the adiabatic regime. See Section 1.5 and the appendix for more details.

In analogy with [Dro19b, Dro19a, DW20], the total number of edge states, signed according to propagation, is the sum over chiralities. From Theorem 3, it is the total Chern number difference. Hence, (1.6) is a form of the bulk edge correspondence, the left-hand-side playing the role of an edge index - see [Hat93, KRSB02, EGS05,


Figure 3. When deforming adiabatically two topological insulators $\mathcal{H}_{0}$ and $\mathcal{H}_{1}$, one must pass discontinuity channels for Chern numbers $\left(s=s_{1}, s_{2}\right)$. These support a signed number of currents equal to the jump of Chern numbers.

ASBVB13, GP13, PSB16, BKR17, Dro19c]. While the interface between $H_{0}$ and $H_{1}$ has width $\delta^{-1}$, the asymmetric transport described above concentrates in finitely many strips of width $\delta^{-1 / 2}$ (corresponding to jumps of local Chern number). This is a much thinner region. This concentration phenomenon - valid only in the adiabatic regime - is not captured by the bulk-edge correspondence.
1.4. General statement. - Theorem 1 will be the consequence of a stronger statement. Let $X$ be a smooth compact manifold of dimension 3 .

Definition 2. - If $H \in C^{\infty}(X, \mathcal{E})$, we say that $H$ has a degeneracy at $x_{0} \in X$ if $H\left(x_{0}\right)$ admits repeated eigenvalues.

We say that this degeneracy is conical if for some $n \in[1, N-1]$ :
(i) $\lambda_{n}\left(H\left(x_{0}\right)\right)=\lambda_{n+1}\left(H\left(x_{0}\right)\right)$ and all other eigenvalues of $H\left(x_{0}\right)$ are simple;
(ii) There exist $\Omega \subset X$ neighborhood of $x_{0}$ in $X, \ell \in C^{\infty}(\Omega, \mathbb{R})$ and $q \in$ $C^{\infty}(\Omega,[0, \infty))$ with a non-degenerate critical value zero at $x_{0}$ such that

$$
\left\{\begin{array}{rl}
\lambda_{n}(H(x)) & =\ell(x)-\sqrt{q(x)}  \tag{1.8}\\
\lambda_{n+1}(H(x)) & =\ell(x)+\sqrt{q(x)}
\end{array}, \quad x \text { near } x_{0} .\right.
$$

For degeneracies of precisely double multiplicity, the mere estimate (1.2) is equivalent to the smooth identity (1.8); see Section 2.1. In other words, Definition 2 corresponds to Definition 1, with the additional requirement (i).

Theorem 3. - When $\operatorname{dim}(X)=3$, the set

$$
\begin{equation*}
\mathbb{M}=\left\{H \in C^{\infty}(X, \mathcal{E}): \text { all degeneracies of } M \text { in } X \text { are conical }\right\} \tag{1.9}
\end{equation*}
$$

is dense and open in $C^{\infty}(X, \varepsilon)$.

According to the von Neumann-Wigner theorem [vNW29], $\mathcal{E} \backslash \mathcal{E}^{*}$ has codimension 3 in $\mathcal{E}$. Since $\operatorname{dim}(X)=3$, the range $H(X)$ of $H$ has Hausdorff dimension at most 3. Thus, generically, $H(X) \cap\left(\mathcal{E} \backslash \mathcal{E}^{*}\right)$ has Hausdorff dimension 0; see Figure 4. This result is closely related to various work about rarity of degenerate eigenvalues in mathematical physics; see e.g. [Col91, Arn95, Tey99]. It is worth mentioning that Theorem 2 applies to non-topologically trivial manifolds - as required in applications to solid state physics. This is the technical part of the proof, see Section 3.4. For related results valid when $X$ is a simply connected regions, we refer to [DP12, §4].

Theorem 3 completes [vNW29]: it shows that the degeneracies of a 3-dimensional family of matrices are conical. In particular, generic elements in $C^{\infty}(X, \mathcal{E})$ have finitely many degeneracies. As an immediate corollary with $X=\mathbb{T}^{3}$ :
Corollary 1. - The degeneracies of Bloch eigenvalues of a generic $\mathbb{Z}^{3}$-invariant Hamiltonian on $\ell^{2}\left(\mathbb{Z}^{3}, \mathbb{C}^{N}\right)$ are all conical.


Figure 4. The range $H(X) \subset \mathcal{E}$ of $H$ has (typical) dimension 3, while $\mathcal{E} \backslash \mathcal{E}^{*}$ has codimension 3. Generically, $H(X)$ and $\mathcal{E} \backslash \mathcal{E}^{*}$ intersect transversely, along a set of dimension 0 .
1.5. Relation with existing work and perspectives. - The present work contrasts with earlier results in tight-binding, quantum graphs, and continuous graphene models [Wal47, Col91, KP07, FW12, AFH ${ }^{+}$18, FLTW18, Lee16]. These papers use the symmetries of the hexagonal lattice to show existence of Dirac points.

The present paper is not symmetry-driven. It is instead topology-driven: conical points arise generically when trying to connect two topologically distinct Hamiltonian, and no other type of degeneracies may form.

When connecting two topologically distinct Hamiltonians, asymmetric currents appear along the interface: the celebrated edge states. Theorems 1 implies that generic edge states of adiabatic systems on $\mathbb{Z}^{2}$ have amplitudes that, after rescaling, evolve according to a universal Dirac-like equation:

$$
\begin{equation*}
\left(D_{t}-\not D\left(x_{2}, D_{x}\right)\right) \beta=0, \quad D_{x}=-i \partial_{x}, \tag{1.10}
\end{equation*}
$$

where $\not D\left(x_{2}, \xi\right)$ is a family of $2 \times 2$ matrices which depends linearly in $x_{2}$ and $\xi$. We refer to the appendix for a formal derivation of (1.10). A full proof would somewhat be transverse to this work; see [FLTW16, Dro19b, ADHY19, DW20, ADH20] for
derivations in slightly different context. See also [FKG03, FK04, Bal19b, Bal19c] for direct work on (1.10).

This Dirac-type propagation should also appear universally in continuous systems - see e.g. [RH08, FLTW16, Dro19b, DW20] for honeycombs. This would require to extend Theorem 3 to differential operators. After some relatively standard reductions, the techniques developed here can treat systems on $L^{2}\left(\mathbb{R}^{2}\right)$ (corresponding to $N=\infty$ ). However they would yield a physically moot genericity result: it would hold within a class much larger than differential operators. We refer to [Col91, Kuc16] for some interesting related conjectures, and formulate our own:

Conjecture 1. - The set
$\left\{V \in C^{\infty}\left(\mathbb{R}^{3} / \mathbb{Z}^{3}\right)\right.$ : all degeneracies of Bloch eigenvalues of $-\Delta_{\mathbb{R}^{3}}+V$ are conical $\}$ is dense and open in $C^{\infty}\left(\mathbb{R}^{3} / \mathbb{Z}^{3}\right)$.
1.6. Organization. - We start with the proof of Theorem 3. In Section 2, we prove that $\mathbb{M}$ is open. This relies on the fact that conical points correspond precisely to critical values zero of the matrix discriminant. In Section 3, we prove that $\mathbb{M}$ is dense. When $N=2$, this boils down to an algebraic identity combined with Sard's theorem. For $N \geqslant 3$, it relies on a reduction to the case $N=2$.

Theorem 1 follows from Theorem 3, as explained in Section 4. The proof of Theorem 2 is independent of the rest of the paper. It relies on arguments from [Dro19a] - see Section 5. In the appendix, we explain the origin of the effective Dirac equation (1.10).

### 1.7. Notations

- Given $N \in \mathbb{N}, \mathcal{E}$ denotes the space of $N \times N$ Hermitian matrices , $\mathcal{E}^{*} \subset \mathcal{E}$ denotes matrices with simple eigenvalues; and $\mathcal{F} \subset \mathcal{E}$ consists of matrices with at most $N-2$ distinct eigenvalues. We provide these spaces with the (Hilbertien) norm $\|A\|^{2}=\operatorname{Tr}_{\mathbb{C}^{N}}\left(A^{2}\right)$.
- Given a smooth compact manifold $X, \mathcal{M}$ is the space $C^{\infty}(X, \varepsilon) ;$ and $\mathbb{M} \subset \mathcal{M}$ consists of elements in $\mathcal{M}$ with only conical degeneracies - see Section 1.4. We fix a Riemannian structure on $X$, with Levi-Civita connection $\nabla$. The space $C^{k}(X, \mathcal{E})$ is the closure of $C^{\infty}(X, \mathcal{E})$ in $C^{0}(X, \mathcal{\varepsilon})$, for the norm

$$
\|H\|_{C^{k}}=\sup \left\{\|H(x)\|+\left\|\nabla^{k} H(x)\right\|: x \in X\right\}, \quad H \in \mathcal{M}=C^{\infty}(X, \mathcal{E}) .
$$

It has a structure of Banach algebra. The space $\mathcal{M}$ inherits a structure of complete metric space, with distance

$$
\begin{equation*}
d(H, \widetilde{H})=\sum_{k=0}^{\infty} 2^{-k} \frac{\|H-\widetilde{H}\|_{C^{k}}}{1+\|H-\widetilde{H}\|_{C^{k}}}, \quad H, \widetilde{H} \in \mathcal{M} . \tag{1.11}
\end{equation*}
$$

- The space $C^{\infty}\left([0,1] \times \mathbb{T}^{2}, \mathcal{E}\right)$ consists of Hermitian-valued smooth functions functions on $(0,1) \times \mathbb{T}^{2}$, whose derivatives extend continuously to $[0,1] \times \mathbb{T}^{2}-$ also provided with the norm (1.11).
- Given $H_{0}, H_{1} \in C^{\infty}\left(\mathbb{T}^{2}, \mathcal{E}\right)$ satisfying (1.1), the space $\mathcal{L} \subset C^{\infty}\left([0,1] \times \mathbb{T}^{2}, \mathcal{E}\right)$ consists of smooth paths connecting $H_{0}$ to $H_{1}$. The space $\mathbb{L} \subset \mathcal{L}$ consists of paths whose $n$-th and $n+1$-th eigenvalues degenerate conically - see Section 1.1.
- The Hausdorff dimension of a set $\mathcal{S}$ is denoted $\operatorname{dim}_{\mathcal{H}}(\mathcal{S})$.
- The Pauli matrices are

$$
\sigma_{0}=\left[\begin{array}{ll}
1 & 0 \\
0 & 1
\end{array}\right], \quad \sigma_{1}=\left[\begin{array}{ll}
0 & 1 \\
1 & 0
\end{array}\right], \quad \sigma_{3}=\left[\begin{array}{cc}
0 & -i \\
i & 0
\end{array}\right], \quad \sigma_{3}=\left[\begin{array}{cc}
1 & 0 \\
0 & -1
\end{array}\right] .
$$

The matrices $\sigma_{1}, \sigma_{2}, \sigma_{3}$ form a basis of the space $\mathcal{E}_{0}$ of traceless Hermitian $2 \times 2$ matrices.

- If $x \in \mathbb{R}^{3}$ and $r>0, \mathbb{B}(x, r)$ is the ball centered at $x$ of radius $r$.


## 2. $\mathbb{M}_{\text {is open }}$

We recall that $\mathcal{M}=C^{\infty}(X, \mathcal{E})$. In this section, we show that the set $\mathbb{M}$ defined in (1.9) is open in $\mathcal{M}$. In Section 2.1 we review the discriminant $D(A)$ of a matrix $A$. This is a quantity depending smoothly on the entries, whose zero set corresponds to matrices with degeneracies.

We then identify conical degeneracies of elements of $\mathcal{M}$ with non-degenerate critical points of $D(H)$. Because of the stability of such points, $\mathbb{M}$ is open in $\mathcal{M}$ - see Section 2.2.
2.1. Discriminant and conical points. - The discriminant of a matrix is the (square of the) Vandermonde determinant of the eigenvalues:

$$
\begin{equation*}
D(A)=\prod_{j \neq k}\left(\lambda_{j}(A)-\lambda_{k}(A)\right)=\prod_{j<k}\left(\lambda_{j}(A)-\lambda_{k}(A)\right)^{2}, \quad A \in \mathcal{E} \tag{2.1}
\end{equation*}
$$

It is a symmetric polynomial in $\lambda_{1}(A), \ldots, \lambda_{N}(A)$. Thus, by the fundamental theorem of linear algebra, it is a polynomial in the quantities $\sum_{j=1}^{m} \lambda_{j}(A)^{m}=\operatorname{Tr}\left[A^{m}\right]-$ see e.g. [Mac15, §I.2]. In particular, $D(A)$ depends smoothly on $A$.

The discriminant detects degenerate eigenvalues: $D(A)=0$ if and only if $A \in \mathcal{E}^{*}$. In fact, it even identifies conical degeneracies.

Lemma 2.1. - $H \in \mathcal{M}$ has a conical degeneracy at $x_{0}$ if and only if $D \circ H$ - i.e., $D \circ H(x)=D(H(x))$ - has a non-degenerate critical value, zero, at $x_{0}$.

Remark 2.1. - No structure - but that of a smooth manifold - is required to define non-degenerate critical points of $u \in C^{\infty}(X, \mathbb{R})$. A Riemannian structure on $X$ allows us to consider the covariant Hessian $\nabla^{2} u$; it is a symmetric endomorphism on $T X$ - see e.g. [Pet16, §2.1]. Non-degenerate critical points correspond to $d u(x)=0$ and $\nabla^{2} u(x)$ non-singular - see e.g. [Pet16, §5.12].

Proof
(1) We assume first that $H$ has a conical degeneracy at $x_{0}$. Let $\lambda_{n}\left(H\left(x_{0}\right)\right)=$ $\lambda_{n+1}\left(H\left(x_{0}\right)\right)$ be the unique degenerate eigenvalue of $H\left(x_{0}\right)$. We write

$$
D \circ H=\left(\lambda_{n+1}(H)-\lambda_{n}(H)\right)^{2} \cdot F, \quad F \stackrel{\text { def }}{=} \prod_{\substack{j<k \\(j, k) \neq(n, n+1)}}\left(\lambda_{j}(H)-\lambda_{k}(H)\right)^{2} .
$$

Using (1.8), $D \circ H=q \cdot F$, where $q \in C^{\infty}(X, \mathbb{R})$ has a non-degenerate critical value zero at $x_{0}$. From general theory, the eigenvalues of Hermitian matrices are Lipschitz in the entries - see [Ser10, Prop. 6.2] - hence a fortiori continuous. Thus $F$ is continuous. Moreover, since all eigenvalues of $H\left(x_{0}\right)$ are simple but $\lambda_{n}\left(H\left(x_{0}\right)\right)=\lambda_{n+1}\left(H\left(x_{0}\right)\right)$, $F\left(x_{0}\right)>0$. We deduce that $D \circ H$ has a non-degenerate critical value zero at $x_{0}$.
(2) Now we assume that $D \circ H$ has the non-degenerate critical value zero at $x_{0}$. Then there exists $\Omega$ neighborhood of $x_{0}$ such that

$$
x \in \Omega \backslash\left\{x_{0}\right\} \Longrightarrow D \circ H(x) \neq 0
$$

In particular, for $x \in \Omega \backslash\left\{x_{0}\right\}$, the eigenvalues $\lambda_{j}(H(x))$ of $H(x)$ are simple - hence smooth functions of $x$.
(3) Since $D \circ H\left(x_{0}\right)=0, H\left(x_{0}\right)$ has at least one degenerate eigenvalue. Define

$$
S=\left\{j \in[1, N-1]: \lambda_{j}\left(H\left(x_{0}\right)\right)=\lambda_{j+1}\left(H\left(x_{0}\right)\right)\right\} .
$$

Since eigenvalues of Hermitian matrices are Lipschitz functions of the entries, there exists $C>0$ such that (after possibly shrinking $\Omega$ ):

$$
\begin{equation*}
x \in \Omega, j \in S \Longrightarrow\left|\lambda_{j}(H(x))-\lambda_{j+1}(H(x))\right| \leqslant C\left\|H(x)-H\left(x_{0}\right)\right\| \tag{2.2}
\end{equation*}
$$

Let $J$ be the cardinal of $S$. From (2.1) and (2.2), we deduce that for some $C^{\prime}>0$,

$$
x \in \Omega \Longrightarrow|D \circ H(x)| \leqslant C^{\prime}\left\|H(x)-H\left(x_{0}\right)\right\|^{2 J}
$$

Since $H$ depends smoothly on $x$ and $D$ has a non-degenerate minimum at $x_{0}$, we deduce that $J \leqslant 1$. This implies that $H\left(x_{0}\right)$ has exactly $N-1$ distinct eigenvalues. Thus, if $n \in[1, N-1]$ is the unique integer such that $\lambda_{n}\left(H\left(x_{0}\right)\right)=\lambda_{n+1}\left(H\left(x_{0}\right)\right)$, then for $j \neq n, n+1, \lambda_{j}(H)$ are smooth in $\Omega$.
(4) Let us fix a contour $\gamma \subset \mathbb{C}$ enclosing $\lambda_{n}\left(H\left(x_{0}\right)\right)=\lambda_{n+1}\left(H\left(x_{0}\right)\right)$ but no other eigenvalue of $H\left(x_{0}\right)$. After possibly shrinking $\Omega$, for $x \in \Omega, \gamma$ enclose $\lambda_{n}(H(x))$ and $\lambda_{n+1}(H(x))$ but no other eigenvalue of $H(x)$. Thus,

$$
\begin{align*}
F_{1}(x) \stackrel{\text { def }}{=} \operatorname{Tr}\left[\int_{\gamma} z(z-H(x))^{-1} \frac{d z}{2 \pi i}\right] & =\lambda_{n}(H(x))+\lambda_{n+1}(H(x))  \tag{2.3}\\
\text { and } \quad F_{2}(x) & \stackrel{\text { def }}{=} \operatorname{Tr}\left[\int_{\gamma} z^{2}(z-H(x))^{-1} \frac{d z}{2 \pi i}\right]
\end{align*}=\lambda_{n}(H(x))^{2}+\lambda_{n+1}(H(x))^{2}
$$

are both smooth functions on $\Omega$. It follows that both

$$
\begin{align*}
& \ell \stackrel{\text { def }}{=} \frac{\lambda_{n}(H)+\lambda_{n+1}(H)}{2}=\frac{F_{1}}{2} \\
& \text { and } \quad q=\frac{\left(\lambda_{n+1}(H)-\lambda_{n}(H)\right)^{2}}{4}=\frac{2 F_{2}-2 F_{1}^{2}}{4} \tag{2.4}
\end{align*}
$$

are smooth functions on $\Omega$.
(5) The equation (2.4) imply that $\lambda_{n}(H)=\ell-\sqrt{q}$ and $\lambda_{n+1}(H)=\ell+\sqrt{q}$. Thus, it remains to show that $q$ has a non-degenerate critical point at $x_{0}$. Again, we write

$$
D \circ H=q \cdot F, \quad F \stackrel{\text { def }}{=} \prod_{\substack{j<k \\(j, k) \neq(n, n+1)}}\left(\lambda_{j}(H)-\lambda_{k}(H)\right)^{2} .
$$

We observe that $F$ is Lipschitz, with $F(0) \neq 0$. Hence, we have

$$
D \circ H(x)=q(x)(1+o(1)) \quad \text { near } x_{0}
$$

this implies

$$
q(x)=D \circ H(x) \cdot(1+o(1))
$$

Since $D \circ H(x)$ has a non-degenerate critical point at $x_{0}$, so does $q$. This completes the proof.
2.2. $\mathbb{M}$ is open. - Here we prove that $\mathbb{M}$ - defined in (1.9) - is open in $\mathcal{M}$. We fix a Riemannian structure on $X$ and consider Hessians of smooth functions on $X$ as symmetric endomorphisms of $T X$ - see Remark 2.1. Define $f: \mathcal{M} \times X \rightarrow \mathbb{R}$ by

$$
\begin{equation*}
f(A, x) \stackrel{\text { def }}{=} \operatorname{Det}\left[\left(\nabla^{2}(D \circ A)\right)(x)\right]^{2}+D \circ A(x) \tag{2.5}
\end{equation*}
$$

Fix $x \in X$ and $H \in \mathbb{M}$. If $H(x) \in \mathcal{E}^{*}$, then $f(H, x) \geqslant D \circ H(x)>0$. If $H(x) \notin \mathcal{E}^{*}$, then $H(x)$ has a conical degeneracy at $x$. Because of Lemma 2.1, $D \circ H$ has a nondegenerate critical point at $x$, thus

$$
f(H, x) \geqslant \operatorname{Det}\left[\left(\nabla^{2}(D \circ H)\right)(x)\right]^{2}>0
$$

We deduce that $f(H, \cdot)$ is positive on $X$; since $X$ is compact, $\inf _{x \in X} f(H, x)>0$.
Since $X$ is compact and $f(A, \cdot)$ depends only on the first two derivatives of $A$, there exists a constant $C$ depending only on $\|H\|_{C^{2}}$ such that

$$
\begin{equation*}
\|B\|_{C^{2}} \leqslant 1 \Longrightarrow|f(H+B, x)-f(H, x)| \leqslant C\|B\|_{C^{2}} \tag{2.6}
\end{equation*}
$$

Since $\inf _{x \in X} f(H, x)>0$, there exists $\varepsilon_{0}>0$ such that whenever $\|B\|_{C^{2}} \leqslant \varepsilon_{0}$, for every $x \in X, f(H+B, x)>0$.

Hence, if $\|B\| \leqslant \varepsilon_{0}$ and $x \in X$, then either:

- $D(H(x)+B(x))>0$, that is $H(x)+B(x) \in \mathcal{E}^{*}$;
- or $D(H(x)+B(x))>0$ and $\operatorname{Det}\left[\left(\nabla^{2} D(H+B)\right)(x)\right]^{2}>0$.

In the latter, $x$ is a non-degenerate critical point of $D(H+B)$. Thus $x$ is a conical degeneracy of $H+B$. This shows that $H+B \in \mathbb{M}$, hence $\mathbb{M}$ is open in $\mathcal{M}$.

## 3. $\mathbb{M}$ is dense

In this section we show that $\mathbb{M}$ is dense in $\mathcal{M}$. When $N=2$, this follows from Sard's theorem and the fact that $D(A)$ is the sum of $3=\operatorname{dim}(X)$ squares depending smoothly on $A$; see Section 3.1.

Two new problems arise for $N \geqslant 3$. Degeneracies can be more intricate: triple eigenvalues or pairs of double eigenvalues may arise. In Section 3.2, we show that
these are too rare to be significant in our problem. This will allow us to focus on $N \times N$ families of matrices with at least $N-1$ distinct eigenvalues.

The other obstacle is more serious: for $N \geqslant 3, D(A)$ is the sum of at least 5 squares - see [Dom11]. Since $5>\operatorname{dim}(X)$, the arguments of Section 3.1 do not naively extend. The key mechanism is that degeneracies of a $N \times N$ family $H \in \mathcal{L}$ with at least $N-1$ distinct eigenvalues reduce locally to those of a $2 \times 2$ family. This enables us to apply the theory of Section 3.1 in simply connected subsets of $X$ - see [DP12, §4] for a related analysis. Passing from these local reductions to a global result on $X$ is the the technical part of the proof, see Section 3.3-3.4. We recall that in solid state physics applications, $X$ is generally not simply connected.
3.1. The case $N=2$. - In this section only, we assume that $N=2$. This considerably simplifies that proof that $\mathbb{M}$ is dense - and it will serve in the general situation.

Proof that $\mathbb{M}$ is dense when $N=2$. - When $N=2$, the Pauli matrices $\sigma_{0}=\operatorname{Id}_{2}, \sigma_{1}$, $\sigma_{2}, \sigma_{3}$ form a basis of $\mathcal{E}$. If $A=\sum_{j=0}^{3} a_{j} \cdot \sigma_{j}$, then

$$
\begin{equation*}
\sigma(A)=a_{0} \pm|a|, \quad D(A)=\|a\|^{2}, \quad \text { where } a=\left[a_{1}, a_{2}, a_{3}\right]^{\top} . \tag{3.1}
\end{equation*}
$$

Let $H \in \mathcal{M}$; we write $H(x)=\sum_{j=0}^{3} h_{j}(x) \cdot \sigma_{j}$. Let $h=\left[h_{1}, h_{2}, h_{3}\right]^{\top}$ and

$$
\begin{aligned}
\mathcal{C} & \stackrel{\text { def }}{=}\left\{t \in \mathbb{R}^{3}: \exists x \in X, h(x)=t \text { and } \operatorname{rk}\left(h^{\prime}(x)\right) \leqslant 2\right\} \\
& =\left\{h(x): x \in X, \operatorname{rk}\left(h^{\prime}(x)\right) \leqslant 2\right\} .
\end{aligned}
$$

According to Sard's theorem, the set $\mathbb{R}^{3} \backslash \mathcal{C}$ is dense in $\mathbb{R}^{3}$ : given $\varepsilon>0$, there exists $b \in \mathbb{R}^{3} \backslash \mathcal{C}$ with $\|b\| \leqslant \varepsilon$; see e.g. [GP74, §1.7]. Set $B=\sum_{j=1}^{3} b_{j} \sigma_{j}$; we claim that all degeneracies of $H-B$ are conical. Indeed from (3.1):

$$
D(H(x)-B)=\|h(x)-b\|^{2}=\sum_{j=1}^{3}\left(h_{j}(x)-b_{j}\right)^{2}
$$

From Lemma 2.1, $H-B$ can have a non-conical degeneracy at a point $x \in X$ only if $h(x)=b$ and $\operatorname{rk}\left(h^{\prime}(x)\right) \leqslant 2$. This is always excluded because $b \notin \mathcal{C}$. Since $\varepsilon$ was arbitrary, we conclude that $\mathbb{M}$ is dense in $\mathcal{M}$ when $N=2$.
3.2. Removing high-multiplicity degeneracies. - We go back to $N \neq 2$. In this section, we explain why we can focus our attention on family of matrices that always have at least $N-1$ distinct eigenvalues.

Lemma 3.1. - The set

$$
\mathcal{F} \stackrel{\text { def }}{=}\{A \in \mathcal{E}: A \text { has at most } N-2 \text { distinct eigenvalues }\}
$$

has Hausdorff dimension at most $N^{2}-6$.
See $[\mathrm{AS} 78, \S 2]$ and $[\mathrm{DE} 99, \S 3]$ for related results - but a different approach. Before giving the proof of Lemma 3.1, we discuss its consequences. We aim to prove that $\mathbb{M}$
is dense in $\mathcal{M}$ : given $H \in \mathcal{M}$ and $\varepsilon>0$, there exists $H_{\varepsilon} \in \mathbb{M}$ such that $d\left(H, H_{\varepsilon}\right) \leqslant 2 \varepsilon$. Since $\operatorname{dim} X=3$ and $\operatorname{dim}_{\mathcal{H}} \mathcal{F} \leqslant N^{2}-6$, the set

$$
\mathcal{S}=\{H(x)-F: x \in X, F \in \mathcal{F}\}
$$

has Hausdorff dimension at most $N^{2}-3$; thus $\mathcal{E} \backslash \mathcal{S}$ has full measure. In particular, there exists $B \in \mathcal{E} \backslash \mathcal{S}$ such that $\|B\| \leqslant \varepsilon$; and $H(x)+B \notin \mathcal{F}$ for every $x \in X$. That is, $H+B \in C^{\infty}(X, \mathcal{E} \backslash \mathcal{F})$.

Thus, to prove that $\mathbb{M}$ is dense in $\mathcal{M}$, we just need to show that for every $H \in$ $C^{\infty}(X, \mathcal{E} \backslash \mathcal{F})$, there exists $H_{\varepsilon} \in \mathbb{M}$ with $d\left(H, H_{\varepsilon}\right) \leqslant \varepsilon$.

Proof of Lemma 3.1
(1) We observe that $\mathcal{F}=\mathcal{F}_{1} \cup \mathcal{F}_{2}$, where

$$
\mathcal{F}_{1}=\{A \in \mathcal{E}: A \text { has a triple eigenvalue }\}, \quad \mathcal{F}_{2}=\mathcal{F} \backslash \mathcal{F}_{1} .
$$

Therefore, it suffices to show that $\mathcal{F}_{1}$ and $\mathcal{F}_{2}$ have Hausdorff dimension at most $N^{2}-6$.
(2) We observe that $\mathcal{F}_{1}=\Phi\left(\mathcal{G}_{1}, \mathbb{R}\right)$, where $\mathcal{G}_{1}$ consists of Hermitian $N \times N$ matrices of rank at most $N-3$; and $\Phi(B, \lambda)=B+\lambda$. We write

$$
\begin{equation*}
\mathcal{G}_{1}=\bigcup_{j=0}^{N-3}\{B \in \mathcal{E}: \operatorname{rk}(B)=j\} \tag{3.2}
\end{equation*}
$$

and we recall that the sets in the RHS of (3.2) are smooth submanifolds of $\mathcal{E}$, of dimension $N^{2}-(N-j)^{2}$ - see e.g. [GP74, §1.4]. Therefore, $\mathcal{G}_{1}$ is a finite union of manifolds of dimensions up to $N^{2}-9$. We deduce that $\operatorname{dim}_{\mathcal{H}} \mathcal{G}_{1}=N^{2}-9$ and $\operatorname{dim}_{\mathcal{H}} \mathcal{F}_{1}=N^{2}-8$.
(3) The set $\mathcal{F}_{2}$ consists of matrices that have two distinct eigenvalues of multiplicity two but no triple eigenvalues. We show that it has Hausdorff dimension at most $N^{2}-6$. For $A_{0} \in \mathcal{F}_{2}$, there exists a unitary $N \times N$ matrix $U$ such that

$$
U^{*} A_{0} U=\left[\begin{array}{ccc}
\lambda_{1} \mathrm{Id}_{2} & 0 & 0 \\
0 & \Lambda & 0 \\
0 & 0 & \lambda_{2} \mathrm{Id}_{2}
\end{array}\right]
$$

where $\lambda_{1} \neq \lambda_{2}$ and $\Lambda$ is a diagonal matrix of size $N-4$, with no diagonal coefficients equal to $\lambda_{1}$ or $\lambda_{2}$. In particular, both

$$
\left[\begin{array}{cc}
\Lambda & 0 \\
0 & \lambda_{2} \mathrm{Id}_{2}
\end{array}\right]-\lambda_{1} \quad \text { and }\left[\begin{array}{cc}
\lambda_{1} \mathrm{Id}_{2} & 0 \\
0 & \Lambda
\end{array}\right]-\lambda_{2}
$$

are invertible $(N-2) \times(N-2)$ matrices. Therefore, there exists a neighborhood $\Omega \subset \mathcal{E}$ of $A_{0}$ such that for any $C \in \Omega$, we can write

$$
U^{*} A U=\left[\begin{array}{ll}
C_{1} & C_{2} \\
C_{2}^{*} & C_{3}
\end{array}\right]=\left[\begin{array}{ll}
D_{3} & D_{2} \\
D_{2}^{*} & D_{1}
\end{array}\right]
$$

where $C_{3}-\lambda_{1}$ and $D_{3}-\lambda_{2}$ are $(N-2) \times(N-2)$ invertible matrices.
(4) If $R_{1}, R_{2}, R_{3}$ are consistently-sized matrices, with $R_{1}$ invertible,

$$
\operatorname{rk}\left(\left[\begin{array}{ll}
R_{1} & R_{2} \\
R_{2}^{*} & R_{3}
\end{array}\right]\right)=\operatorname{rk}\left(R_{1}\right)+\operatorname{rk}\left(R_{3}-R_{2}^{*} R_{1}^{-1} R_{2}\right)
$$

This can be seen for instance from Schur's complement formula:

$$
\left[\begin{array}{ll}
R_{1} & R_{2}  \tag{3.3}\\
R_{2}^{*} & R_{3}
\end{array}\right]\left[\begin{array}{cc}
\operatorname{Id}_{N-2} & -R_{1}^{-1} R_{2} \\
0 & \operatorname{Id}_{2}
\end{array}\right]=\left[\begin{array}{cc}
R_{1} & 0 \\
R_{2}^{*} & R_{3}-R_{2}^{*} R_{1}^{-1} R_{2}
\end{array}\right] .
$$

Let $\Omega_{1}, \Omega_{2} \subset \mathbb{R}$ be sufficiently small disjoint neighborhood of $\lambda_{1}, \lambda_{2}$ such that if $\mu_{1} \in \Omega_{1}$ and $\mu_{2} \in \Omega_{2}$,

$$
\Phi_{\mu_{1}, \mu_{2}}(A)=\left(C_{1}-\mu_{1}-C_{2}^{*}\left(C_{3}-\mu_{1}\right)^{-1} C_{2}, D_{1}-\mu_{2}-D_{2}^{*}\left(D_{3}-\mu_{2}\right)^{-1} D_{2}\right),
$$

from $\Omega$ to pairs of $2 \times 2$ Hermitian matrices, is well-defined. By (3.3), $\Phi_{\mu_{1}, \mu_{2}}(A)=(0,0)$ if and only if $A-\mu_{1}$ and $A-\mu_{2}$ are of rank $N-2$; equivalently, if and only if $\mu_{1}$ and $\mu_{2}$ are two double eigenvalues of $A$.
(5) The map $\Phi_{\mu_{1}, \mu_{2}}$ is a local submersion at $A_{0}$. Indeed, we have

$$
d \Phi_{\mu_{1}, \mu_{2}}\left(A_{0}\right) \cdot U\left[\begin{array}{ccc}
\epsilon_{1} & 0 & 0 \\
0 & 0 & 0 \\
0 & 0 & \epsilon_{2}
\end{array}\right] U^{*}=\left(\epsilon_{1}, \epsilon_{2}\right)
$$

We note that $\Phi_{\mu_{1}, \mu_{2}}$ has range in pairs of $2 \times 2$ Hermitian matrices, which has dimension 8 . Thus, by the local submersion theorem [GP74, §4], $\Phi_{\mu_{1}, \mu_{2}}^{-1}(0,0)$ is a submanifold of $\mathcal{E}$ of dimension $N^{2}-8$.

Using continuity of eigenvalues, after potentially shrinking $\Omega$, we have

$$
\mathcal{F}_{2} \cap \Omega=\underset{\left(\mu_{1}, \mu_{2}\right) \in \Omega_{1} \times \Omega_{2}}{ } \Phi_{\mu_{1}, \mu_{2}}^{-1}(0,0) .
$$

Since $\Omega_{1} \times \Omega_{2}$ has dimension 2 , $\operatorname{dim}_{\mathcal{H}}\left(\mathcal{F}_{2} \cap \Omega\right) \leqslant N^{2}-6$. Since $\Omega \subset \mathcal{E}$ is a neighborhood of an arbitrary element $A_{0} \in \mathcal{F}_{2}, \mathcal{F}_{2}$ is a countable union of sets of dimension at most $N^{2}-6$, thus it has dimension at most $N^{2}-6$.
3.3. Removing bad points: preparatory lemmas. - Because of Section 3.2, we focus (without loss of generalities) on $H \in C^{\infty}(X, \mathcal{E} \backslash \mathcal{F}): H$ has, at all points of $X$, at least $N-1$ distinct eigenvalues. We will show in Section 3.4 that $H$ is arbitrarily close to $\mathbb{M}$.

A naive generalization of Section 3.1 to $N \geqslant 3$ requires to write $D(A)$ as a sum of three squares depending smoothly on $A \in \mathcal{E}$ - see (3.1). This is not possible: according to [Dom11],‘ at least 5 squares are necessary; see also [Ily92, Lax98, Par02, Dom11]. In Section 3.4, we will get around by writing $D(H)$ locally - instead of globally - as a sum of 3 squares. The present section lays out preparatory lemmas.

Fix $x_{\star} \in X$. According to the assumption, there exists $n_{\star} \in[1, N-1]$ such that

$$
\lambda_{1}\left(H\left(x_{\star}\right)\right)<\cdots<\lambda_{n_{\star}}\left(H\left(x_{\star}\right)\right) \leqslant \lambda_{n_{\star}+1}\left(H\left(x_{\star}\right)\right)<\cdots<\lambda_{N}\left(H\left(x_{\star}\right)\right) .
$$

Since eigenvalues are continuous functions of the entries, there exists an open neighborhood $X_{\star} \subset X$ of $x_{\star}$ such that
(3.4) $x \in X_{\star} \Longrightarrow \lambda_{1}(H(x))<\cdots<\lambda_{n_{\star}}(H(x)) \leqslant \lambda_{n_{\star}+1}(H(x))<\cdots<\lambda_{N}(H(x))$.

After potentially shrinking $X_{\star}$, there exists a ball $B\left(0,2 r_{\star}\right) \subset \mathbb{R}^{3}$, and a smooth diffeomorphism $\phi_{\star}: B\left(0,2 r_{\star}\right) \rightarrow X_{\star}$ with $\phi_{\star}(0)=x_{\star}$. We set $Y_{\star}=\phi\left(B\left(0, r_{\star}\right)\right) \subset X_{\star}$.

We observe that $x_{\star} \in Y_{\star}$. Thus, the collection of open sets $\left\{Y_{\star}\right\}_{x_{\star} \in X}$ covers $X$ and we can pass to a finite collection, associated to points $x_{1}, \ldots, x_{P}$.

Lemma 3.2. - There exists $\delta_{0} \in(0,1)$ such that for all $B \in \mathcal{M}$ with $\|H-B\|_{C^{0}} \leqslant \delta_{0}$, for every $p \in[1, P]$,

$$
x \in X_{p} \Longrightarrow \lambda_{1}(B(x))<\cdots<\lambda_{n_{p}}(B(x)) \leqslant \lambda_{n_{p}+1}(B(x))<\cdots<\lambda_{N}(B(x)) .
$$

This result is a direct consequence of (3.4) with continuity of eigenvalues in the entries of the matrix - [Ser10, Prop. 6.2].

Given $A \in \mathcal{M}$, we say that $x \in X$ is a bad point of $A$ if $A$ has a non-conical degeneracy at $x$. We let $\mathscr{B}(A)$ be the set of bad points of $A$; in particular, $A \in \mathbb{M}$ if and only if $\mathscr{B}(A)=\varnothing$. Bad points are stable:

Lemma 3.3. - Let $A \in \mathcal{M}$ and $Z \subset X$ be an open set such that $\mathscr{B}(A) \subset Z$. Then there exists $\eta_{0} \in(0,1)$ such that for all $B \in \mathcal{M}$ with $\|B\|_{C^{2}} \leqslant \eta_{0}, \mathscr{B}(A+B) \subset Z$.

Proof. - Recall (2.5) and (2.6): there exists $C>0$ (depending on $\|A\|_{C^{2}}$ ) such that

$$
\|B\|_{C^{2}} \leqslant 1 \Longrightarrow|f(A+B, x)-f(A, x)| \leqslant C\|B\|_{C^{2}}
$$

$$
\begin{equation*}
\text { where } f(A, x) \stackrel{\text { def }}{=} \operatorname{Det}\left[\left(\nabla^{2}(D \circ A)\right)(x)\right]^{2}+D \circ A(x) \tag{3.5}
\end{equation*}
$$

Moreover, $f(A+B, x)=0$ if and only if $x \in \mathscr{B}(A+B)$.
On the compact set $X \backslash Z, f(A, \cdot)>0$. From (3.5), if $\|B\|_{C^{2}}$ is sufficiently small, $f(A+B, \cdot)>0$ on $X \backslash Z$. Thus $\mathscr{B}(A+B) \subset Z$. This completes the proof.
3.4. Proof of Theorem 3. - We refer to Figure 5 for a step-by-step pictorial explanation of the proof.

## Proof that $\mathbb{M}$ is dense in $\mathcal{M}$

(1) As explained in Section 3.2, to prove density of $\mathbb{M}$ in $\mathcal{M}$, it suffices to prove density of $\mathbb{M}$ in $C^{\infty}(X, \mathcal{E} \backslash \mathcal{F})$. Let $H \in C^{\infty}(X, \mathcal{E} \backslash \mathcal{F})$. Fix $0<\varepsilon<\delta_{0} / 4$, where $\delta_{0}$ is given by Lemma 3.2. For each $p \in[0, P]$, we construct recursively $H_{p} \in \mathcal{M}$ such that

$$
d\left(H, H_{p}\right) \leqslant\left(1-2^{-p}\right) \varepsilon ; \quad \text { and } \quad \mathscr{B}\left(H_{p}\right) \subset Z_{p}, \quad Z_{p} \stackrel{\text { def }}{=} Y_{p+1} \cup \cdots \cup Y_{P}
$$

In particular, $H_{P}$ will satisfy $d\left(H, H_{P}\right) \leqslant \varepsilon$ and $\mathscr{B}\left(H_{P}\right)=\varnothing$.
For $p=0$, we simply take $H_{0}=H$. For $p \geqslant 1$, we proceed by induction: we assume that $H_{p-1}$ is constructed and we want to construct $H_{p}$.
(2) For $x \in X_{p}$, let $\mathcal{V}(x)$ be the eigenspace of $H_{p-1}(x)$ associated to the eigenvalues $\lambda_{n_{p}}\left(H_{p-1}(x)\right)$ and $\lambda_{n_{p}+1}\left(H_{p-1}(x)\right)$. Since $d\left(H, H_{p-1}\right) \leqslant \varepsilon,\left\|H-H_{p-1}\right\|_{C^{0}} \leqslant \delta_{0}$. Thus Lemma 3.2 implies that for every $x \in X_{p}$,

$$
\begin{equation*}
\lambda_{1}\left(H_{p-1}(x)\right)<\cdots<\lambda_{n_{p}}\left(H_{p-1}(x)\right) \leqslant \lambda_{n_{p}+1}\left(H_{p-1}(x)\right)<\cdots<\lambda_{N}\left(H_{p-1}(x)\right) . \tag{3.6}
\end{equation*}
$$

Because of (3.6), $\mathcal{V}(x)$ induces a rank-two vector bundle over $X_{p}$; and so does $\mathcal{V}(x)^{\perp}$. Since $X_{p}$ is diffeomorphic to a ball in $\mathbb{R}^{3}, \mathcal{V}$ and $\mathcal{V}^{\perp}$ are trivial vector bundles - see

(d) We get new systems $H_{1}, H_{2}, H_{3}$, recursively constructed, arbitrarily close to $H$, with no bad points in $Y_{1}, Y_{1} \cup Y_{2} \backslash Y_{3}, Y_{1} \cup Y_{2} \cup Y_{3}$, respectively. Since $Y_{1} \cup Y_{2} \cup Y_{3}$ cover $X$, $H_{3}$ is in $\mathbb{M}$ and is arbitrarily close to $H$.
e.g. [Moo01, §1.3]. Therefore, they both admit unitary frames. This means that there exists $U \in C^{\infty}\left(X_{p}, U(N)\right)$ such that for all $x \in X_{p}$,

$$
H_{p-1}(x)=U(x)\left[\begin{array}{cc}
J(x) & 0  \tag{3.7}\\
0 & J(x)^{\perp}
\end{array}\right] U(x)^{*}
$$

where

- $J(x)$ is a $2 \times 2$ Hermitian matrix depending smoothly on $x \in X_{p}$, with eigenvalues $\lambda_{n_{p}}\left(H_{p-1}(x)\right)$ and $\lambda_{n_{p}+1}\left(H_{p-1}(x)\right)$;
$-J(x)^{\perp}$ is a $(N-2) \times(N-2)$ Hermitian matrix depending smoothly on $x \in X_{p}$, with simple eigenvalues $\lambda_{j}\left(H_{p-1}(x)\right), j \notin\left\{n_{p}, n_{p}+1\right\}$.
(3) Let $\chi \in C^{\infty}(X, \mathbb{R})$ be equal to 1 on a neighborhood of $Y_{p}$, with support contained in $X_{p}$. Let $B$ be a Hermitian $2 \times 2$ matrix and define

$$
H_{p}(x) \stackrel{\text { def }}{=} H_{p-1}(x)+\chi(x)^{2} \cdot U(x)\left[\begin{array}{cc}
B & 0 \\
0 & 0
\end{array}\right] U(x)^{*}
$$

We note that $H_{p} \in \mathcal{M}: \chi=0$ when $U$ is not well-defined. As $C^{k}(X, \mathbb{R})$ is an algebra,

$$
\left\|H_{p}-H_{p-1}\right\|_{C^{k}} \leqslant \alpha_{k}\|B\|, \quad \alpha_{k} \stackrel{\text { def }}{=} C_{k}\|\chi U\|_{C^{k}}\left\|\chi U^{*}\right\|_{C^{k}}
$$

Using that $s \mapsto s(1+s)^{-1}$ increases on $[0, \infty)$,

$$
\begin{equation*}
d\left(H_{p}, H_{p-1}\right)=\sum_{k=0}^{\infty} 2^{-k} \frac{\left\|H_{p}-H_{p-1}\right\|_{C^{k}}}{1+\left\|H_{p}-H_{p-1}\right\|_{C^{k}}} \leqslant \sum_{k=0}^{\infty} 2^{-k} \frac{\alpha_{k}\|B\|}{1+\alpha_{k}\|B\|} \tag{3.8}
\end{equation*}
$$

We split the sum in the RHS in two parts, depending whether $\alpha_{k}$ is larger than $\|B\|^{-1 / 2}$. Since $s(1+s)^{-1} \leqslant \min (1, s)$, we deduce that

$$
\sum_{\alpha_{k} \leqslant\|B\|^{-1 / 2}} 2^{-k} \frac{\alpha_{k}\|B\|}{1+\alpha_{k}\|B\|} \leqslant 2\|B\|^{1 / 2}, \quad \sum_{\alpha_{k}>\|B\|^{-1 / 2}} 2^{-k} \frac{\alpha_{k}\|B\|}{1+\alpha_{k}\|B\|} \leqslant 2^{-k_{B}+1}
$$

where $k_{B}$ is the smallest integer such that $\alpha_{k}>\|B\|^{-1 / 2}$ (with $k_{B}=\infty$ if no such integer exist). In particular, $k_{B} \rightarrow \infty$ as $\|B\| \rightarrow 0$. Going back to (3.8), we deduce that

$$
\begin{equation*}
d\left(H_{p}, H_{p-1}\right) \leqslant 2\left(\|B\|^{1 / 2}+2^{-k_{B}}\right) \longrightarrow 0 \quad \text { as } \quad\|B\| \longrightarrow 0 \tag{3.9}
\end{equation*}
$$

(4) Let $\eta_{0}$ be associated to $H_{p-1}$ and $Z_{p-1}$ by Lemma 3.3. Thanks to Section 3.1 and (3.9), we can find a Hermitian $2 \times 2$ matrix $B$ with the two following conditions:

- All degeneracies of $J(x)+B$ in $X_{p}$ are conical;
$-d\left(H_{p}, H_{p-1}\right) \leqslant \min \left(2^{-p} \varepsilon, \eta_{0} / 8\right)$.
The recursion assumption $d\left(H, H_{p-1}\right) \leqslant\left(1-2^{-p-1}\right) \varepsilon$ and $d\left(H_{p}, H_{p-1}\right) \leqslant 2^{-p} \varepsilon$ yield $d\left(H, H_{p}\right) \leqslant\left(1-2^{-p}\right) \varepsilon$. Moreover, $d\left(H_{p}, H_{p-1}\right) \leqslant \eta_{0} / 8$ implies $\left\|H_{p}-H_{p-1}\right\|_{C^{2}} \leqslant \eta_{0}$. From Lemma 3.3 and the recursion assumption $\mathscr{B}\left(H_{p-1}\right) \subset Z_{p-1}, \mathscr{B}\left(H_{p}\right) \subset Z_{p-1}$.
(5) To complete the recursion, it remains to show that $\mathscr{B}\left(H_{p}\right) \subset Z_{p}$; equivalently, that $H_{p}$ has no bad degeneracies in $Y_{p}$. When $\chi(x)=1$ (i.e., on a neighborhood of $Y_{p}$ ),

$$
H_{p}(x)=U(x)\left[\begin{array}{cc}
J(x)+B & 0  \tag{3.10}\\
0 & J(x)^{\perp}
\end{array}\right] U(x)^{*}
$$

Using (3.7), the identity (3.10) implies that when $\chi(x)=1$, the eigenvalues of $H_{p}(x)$ are: $\lambda_{j}\left(H_{p-1}(x)\right)$ for $j \neq n_{p}, n_{p}+1$; and $\lambda_{j}(J(x)+B), j=1,2$.

From (3.6), the only possible degeneracies of $H_{p}$ in $\{\chi=1\}$ arise from $\lambda_{1}(J+B)$ and $\lambda_{2}(J+B)$. By definition of $B$, all such degeneracies are conical. Since $Y_{p} \subset$ $\{\chi=1\}$, we get $\mathscr{B}\left(H_{p}\right) \cap Y_{p}=\varnothing$. This completes the recursion and the proof of Theorem 3.

## 4. Proof of Theorem 1

Proof that $\mathbb{L}$ is open in $\mathcal{L}$. - The proof is similar to Section 2. Fix $H \in \mathbb{L}$; let $\left\{\zeta_{1}, \ldots, \zeta_{J}\right\}$ be the (finite) set of points of $[0,1] \times \mathbb{T}^{2}$, such that $\lambda_{n}(H)$ and $\lambda_{n+1}(H)$ degenerate.

For each $j \in[1, J]$, let $\gamma_{j}$ be a contour enclosing $\lambda_{n}\left(H\left(\zeta_{j}\right)\right)=\lambda_{n+1}\left(H\left(\zeta_{j}\right)\right)$, but no other eigenvalue of $H\left(\zeta_{j}\right)$. Using continuity of eigenvalues, there exist $\varepsilon_{0}$ and $r_{0}>0$ such that for $B \in \mathcal{L}$ with $\|B\|_{C^{2}} \leqslant \varepsilon_{0}$ and $\zeta \in \mathbb{B}\left(\zeta_{j}, r_{0}\right), \gamma_{j}$ encloses $\lambda_{n}(H(\zeta)+B(\zeta))$ and $\lambda_{n+1}(H(\zeta)+B(\zeta))$ but no other eigenvalues of $H(\zeta)+B(\zeta)$.

Without loss of generality, the balls $\mathbb{B}\left(\zeta_{j}, r_{0}\right)$ are disjoints. For $\zeta \in \mathbb{B}\left(\zeta_{j}, r_{0}\right)$, introduce, similarly to (2.3),

$$
\begin{aligned}
& G_{1}(\zeta, B) \stackrel{\text { def }}{=} \operatorname{Tr}\left[\int_{\gamma_{j}} z(z-H(\zeta)-B(\zeta))^{-1} \frac{d z}{2 \pi i}\right]=\sum_{j=n}^{n+1} \lambda_{j}(H(\zeta)+B(\zeta)), \\
& G_{2}(\zeta, B) \stackrel{\text { def }}{=} \operatorname{Tr}\left[\int_{\gamma_{j}} z^{2}(z-H(\zeta)-B(\zeta))^{-1} \frac{d z}{2 \pi i}\right]=\sum_{j=n}^{n+1} \lambda_{j}(H(\zeta)+B(\zeta))^{2}, \\
& G(\zeta, B) \stackrel{\text { def }}{=} 2 G_{2}(\zeta, B)-G_{1}(\zeta, B)^{2}=\left(\lambda_{n+1}(H(\zeta)+B(\zeta))-\lambda_{n}(H(\zeta)+B(\zeta))\right)^{2} .
\end{aligned}
$$

We note that $G\left(\zeta_{j}, 0\right)=0$ hence $\nabla_{\zeta}^{2} G\left(\zeta_{j}, 0\right)>0$, because $\lambda_{n}(H)$ and $\lambda_{n+1}(H)$ may only degenerate conically. The identity $G=2 G_{2}-G_{1}^{2}$ and the Cauchy representation of $G_{1}$ and $G_{2}$ imply that for some $C>0$ and all $\zeta \in \Omega$,

$$
\left|\nabla_{\zeta}^{2} G(\zeta, B)-\nabla_{\zeta}^{2} G(\zeta, 0)\right| \leqslant C\|B\|_{C^{2}}
$$

Therefore, after possibly shrinking $\varepsilon_{0}$ and $\Omega$,

$$
\|B\|_{C^{2}} \leqslant \varepsilon_{0}, \zeta \in \Omega \Longrightarrow \nabla_{\zeta}^{2} G(\zeta, B)>0 .
$$

Thus, if $\lambda_{n}(H(\zeta)+B(\zeta))=\lambda_{n+1}(H(\zeta)+B(\zeta))$ for $\zeta \in \Omega$, then this degeneracy is conical. Finally, after shrinking $\varepsilon_{0}, \lambda_{n}(H+B)$ and $\lambda_{n+1}(H+B)$ cannot degenerate outside $\Omega$. This shows that $H+B \in \mathbb{L}$ : $\mathbb{L}$ is open in $\mathcal{L}$.

Proof that $\mathbb{L}$ is dense in $\mathcal{L}$
(1) We show that $\mathbb{L}$ is dense in $\mathcal{L}$. Since eigenvalues are Lipschitz functions of the matrix entries, we deduce from (1.1) that there exists $\eta_{0} \in(0,1)$ such that for every $T \in \mathcal{E}$,

$$
\|T\| \leqslant 4 \eta_{0}, \quad \xi \in \mathbb{T}^{2} \Longrightarrow\left\{\begin{array}{l}
\lambda_{n+1}\left(H_{0}(\xi)+T\right)>\lambda_{n}\left(H_{0}(\xi)+T\right)  \tag{4.1}\\
\lambda_{n+1}\left(H_{1}(\xi)+T\right)>\lambda_{n}\left(H_{1}(\xi)+T\right)
\end{array} .\right.
$$

(2) Let $H \in \mathcal{L}: H$ is smooth on $(0,1) \times \mathbb{T}^{2}$, with bounded derivatives; and connects $H_{0}$ to $H_{1}$. Seeley's operator [See64] extends $H$ as an element of $C_{0}^{\infty}\left((-\pi, \pi) \times \mathbb{T}^{2}, \mathcal{E}\right)$, thus as an element of $C^{\infty}\left(\mathbb{T}^{3}, \mathcal{E}\right)$ (still denoted $H$ ).

Let $\chi_{0}, \chi_{1} \in C_{0}^{\infty}\left(\mathbb{T}^{1},[0,1]\right)$ with $\chi_{0}(0)=\chi_{1}(1)=1$ and
$\operatorname{supp}\left(\chi_{0}\right) \subset\left(-\delta_{0}, \delta_{0}\right) /(2 \pi \mathbb{Z}), \quad \operatorname{supp}\left(\chi_{0}\right) \subset\left(1-r_{0}, 1+r_{0}\right) /(2 \pi \mathbb{Z}), \quad r_{0} \stackrel{\text { def }}{=} \frac{\eta_{0}}{1+\|H\|_{C^{1}}}$.
For $\mathcal{H} \in C^{\infty}\left(\mathbb{T}^{3}, \mathcal{E}\right)$, we introduce

$$
\begin{equation*}
\mathbb{H}(s, \xi)=\mathcal{H}(s, \xi)+\chi_{0}(s)\left(H_{0}(\xi)-\mathcal{H}(0, \xi)\right)+\chi_{1}(s)\left(H_{1}(\xi)-\mathcal{H}(1, \xi)\right) \tag{4.2}
\end{equation*}
$$

We observe that $\mathbb{H}$ restricts to $[0,1] \times \mathbb{T}^{2}$ as an element of $\mathcal{L}$ : it varies smoothly with $(s, \xi)$ and connects $H_{0}$ to $H_{1}$.
(3) Fix $\varepsilon>0$. Using (4.2) and that $C^{k}\left(\mathbb{T}^{3}, \mathcal{E}\right)$ is an algebra, we have $\|H-\mathbb{H}\|_{C^{k}} \leqslant$ $C_{k}\|H-\mathcal{H}\|_{C^{k}}$ for some $C_{k}>0$. As in Step 3 in Section 3.4 there exists $\eta_{1} \in\left(0, \eta_{0}\right)$ with

$$
d(H, \mathcal{H}) \leqslant \eta_{1} \Longrightarrow d(H, \mathbb{H}) \leqslant \varepsilon .
$$

We now demand that $\mathcal{H} \in \mathbb{M}$, and $d(H, \mathcal{H}) \leqslant \eta_{1}$; such $\mathcal{H}$ exist by Theorem 3. Under these conditions, $\mathbb{H}$ defined by (4.2) satisfies $d(H, \mathbb{H}) \leqslant \varepsilon$; we claim that $\lambda_{n}(\mathbb{H})$ and $\lambda_{n+1}(\mathbb{H})$ can only degenerate conically in $[0,1] \times \mathbb{T}^{2}$.
(4) For $(s, \xi) \in\left(r_{0}, 1-r_{0}\right) \times \mathbb{T}^{2}$, we have $\mathbb{H}(s, \xi)=\mathcal{H}(s, \xi)$. Since $\mathcal{H} \in \mathbb{M}$, we deduce that $\lambda_{n}(\mathbb{H})$ and $\lambda_{n+1}(\mathbb{H})$ can only degenerate conically in $\left(r_{0}, 1-r_{0}\right) \times \mathbb{T}^{2}$.

For $(s, \xi) \in\left[0, r_{0}\right] \times \mathbb{T}^{2}$, we have

$$
\begin{aligned}
\left\|\mathbb{H}(s, \xi)-H_{0}(\xi)\right\| & \leqslant\left\|\mathcal{H}(s, \xi)-H_{0}\right\|+\left\|H_{0}(\xi)-\mathcal{H}(0, \xi)\right\| \\
& \leqslant\|\mathcal{H}(s, \xi)-\mathcal{H}(0, \xi)\|+2\left\|H_{0}(\xi)-\mathcal{H}(0, \xi)\right\| \\
& \leqslant r_{0}\|\mathcal{H}\|_{C^{1}}+2 \delta \leqslant r_{0}\left(\delta+\|H\|_{C^{1}}\right)+2 \delta \leqslant 3 \eta_{0} .
\end{aligned}
$$

In the last line, we used the definition of $r_{0}$ and the inequality $\delta<\eta_{0}<1$. Thanks to (4.1), we deduce that $\lambda_{n}(\mathbb{H})$ and $\lambda_{n+1}(\mathbb{H})$ cannot cross in $\left[0, r_{0}\right] \times \mathbb{T}^{2}$. A similar argument shows that they cannot cross in $\left[1-r_{0}, 1\right] \times \mathbb{T}^{2}$.

Hence, the restriction of $\mathbb{H}$ to $[0,1] \times \mathbb{T}^{2}$ is in $\mathbb{L}$; and $d(H, \mathbb{H}) \leqslant \varepsilon$. Since $\varepsilon$ was arbitrary, $\mathbb{L}$ is dense in $\mathcal{L}$. This completes the proof of Theorem 1 .

## 5. Chern number difference

## Proof of Theorem 2

(1) We start with a few notations and definitions. Let $H \in \mathbb{L}$. Let $\mathcal{R}$ be the set of points $\zeta=(s, \xi) \in[0,1] \times \mathbb{T}^{2}$ such that $\lambda_{n}(H(\zeta))<\lambda_{n+1}(H(\zeta))$. For $\zeta \in \mathcal{R}$, we can represent the projector $\Pi_{n}(\zeta)$ to the first $n$ eigenspaces of $H(\zeta)$ as a Cauchy integral:

$$
\begin{equation*}
\Pi_{n}(\zeta)=\frac{1}{2 \pi i} \oint_{\gamma_{n}(\zeta)}(z-H(\zeta))^{-1} d z \tag{5.1}
\end{equation*}
$$

where $\gamma_{n}(\zeta) \subset \mathbb{C}$ encloses $\lambda_{1}(H(\zeta)), \ldots, \lambda_{n}(H(\zeta))$ but no other eigenvalue of $H(\zeta)$. If $\{s\} \times \mathbb{T}^{2} \subset \mathcal{R}$, then $\Pi_{n}(s, \cdot)$ induces a vector bundle over $\mathbb{T}^{2}$ : the fiber at $\xi \in \mathbb{T}^{2}$ is

Range $\left(\Pi_{n}(s, \xi)\right)$. We let $\mathcal{S}$ be the set of $s \in[0,1]$ such that $\mathscr{V}_{s}$ is not well-defined equivalently, $\mathcal{S}=\left\{s \in[0,1]: \exists \xi \in \mathbb{T}^{2},(s, \xi) \notin \mathcal{R}\right\}$.

For $\zeta \in \mathcal{R}$, we define

$$
\begin{equation*}
B_{n}(\zeta)=\operatorname{Tr}_{\mathbb{C}^{N}}\left(\Pi_{n}(\zeta)\left[\partial_{\xi_{1}} \Pi_{n}(\zeta), \partial_{\xi_{2}} \Pi_{n}(\zeta)\right) .\right. \tag{5.2}
\end{equation*}
$$

This is a smoothly varying function on $\mathcal{R}$, that interprets as the Berry curvature. In particular, $B_{n}(\xi) d \xi$ is a two-form; and $B_{n}(\xi)$ is additive: if $\zeta \in \mathcal{R}$ and $\lambda_{n-1}(H(\zeta))<$ $\lambda_{n}(H(\zeta))$, then $B_{n}(\zeta)=B_{n-1}(\zeta)+b(\zeta)$, where:

- $B_{n-1}(\zeta)$ is associated with the projector $\Pi_{n-1}(\zeta)$ to the first $n-1$ eigenspaces of $H(\zeta)$ - see (5.1), (5.2) with $n$ replaced by $n-1$;
$-b(\zeta)$ is associated to the rank-one projector $\pi(\zeta)$ to $\operatorname{ker}\left(\lambda_{n}(H(\zeta))-H(\zeta)\right)$ :

$$
b(\zeta)=\operatorname{Tr}_{\mathbb{C}^{N}}\left(\pi(\zeta)\left[\partial_{\xi_{1}} \pi(\zeta), \partial_{\xi_{2}} \pi(\zeta)\right]\right) .
$$

For $s \in[0,1] \backslash \mathcal{S}$, the Chern number of $\mathscr{V}_{s}$ is the integer

$$
c_{1}\left(\mathscr{V}_{s}\right)=\frac{1}{2 \pi i} \int_{\mathbb{T}^{2}} B_{n}(s, \xi) d \xi
$$

In Step 6 , we will use the space of $2 \times 2$ traceless Hermitian matrices $\mathcal{E}_{0}$. This space is equipped with the Hermitian inner product $\left\langle T_{1}, T_{2}\right\rangle=\operatorname{Tr}\left(T_{1} T_{2}\right)$; the Pauli matrices $\sigma_{1}, \sigma_{2}, \sigma_{3}$ form an orthonormal basis. If $\widetilde{\sigma}_{1}, \widetilde{\sigma}_{2}, \widetilde{\sigma}_{3}$ is another orthonormal basis, then there exists $U \in S U(2)$ (unique up to multiplication by $\pm \operatorname{Id}_{2}$ ) and $\epsilon \in\{ \pm 1\}$ such that

$$
\begin{equation*}
\widetilde{\sigma}_{k}=\epsilon \cdot U \sigma_{k} U^{*}, \quad 1 \leqslant k \leqslant 3 \tag{5.3}
\end{equation*}
$$

This is precisely the content of the isomorphism between $S U(2) /\left\{ \pm \mathrm{Id}_{2}\right\}$ and $S O(3)$; see e.g. [Sin05, §4.2]. The number $\epsilon \in\{ \pm 1\}$ reads as the determinant of the (orthogonal) matrix of the basis ( $\left.\widetilde{\sigma}_{1}, \widetilde{\sigma}_{2}, \widetilde{\sigma}_{3}\right)$ in the basis $\left(\sigma_{1}, \sigma_{2}, \sigma_{3}\right)$.
(2) Since $H \in \mathbb{L}$, the sets $[0,1] \backslash \mathcal{R}$ and $\mathcal{S}$ are finite. The map $s \mapsto c_{1}\left(\mathscr{V}_{s}\right)$ is welldefined on $[0,1] \backslash \mathcal{S}$. Since it is integer-valued, it is locally constant on each sub-interval of $[0,1] \backslash \mathcal{S}$. We deduce that

$$
\begin{align*}
c_{1}\left(\mathscr{V}_{1}\right)-c_{1}\left(\mathscr{V}_{0}\right) & =\lim _{\delta \rightarrow 0^{+}} \sum_{s_{\star} \in \mathcal{S}} c_{1}\left(\mathscr{V}_{s_{\star}+\delta}\right)-c_{1}\left(\mathscr{V}_{s_{\star}-\delta}\right) \\
& =\frac{1}{2 \pi i} \sum_{s_{\star} \in \mathcal{S}} \lim _{\delta \rightarrow 0^{+}} \int_{\mathbb{T}^{2}}\left(B_{n}\left(s_{\star}+\delta, \xi\right)-B_{n}\left(s_{\star}-\delta, \xi\right)\right) d \xi . \tag{5.4}
\end{align*}
$$

It remains to compute each individual summand in the RHS of (5.4). For that, we use the techniques developed in [Dro19a, §2] - and we refer to that paper for full details.
(3) Fix $s_{\star} \in \mathcal{S}$; let $\mathcal{Z}$ be the set of points $\xi \in \mathbb{T}^{2}$ such that $\left(s_{\star}, \xi\right) \notin \mathcal{R}$. Using that $B(\zeta)$ depends smoothly on $\zeta \in \mathcal{R}$, we deduce that for $r$ sufficiently small,

$$
\begin{align*}
\int_{\mathbb{T}^{2}}\left(B_{n}\left(s_{\star}+\delta, \xi\right)-\right. & \left.B_{n}\left(s_{\star}-\delta, \xi\right)\right) d \xi  \tag{5.5}\\
& =\sum_{\xi_{\star} \in Z} \int_{\left|\xi-\xi_{\star}\right| \leqslant r}\left(B_{n}\left(s_{\star}+\delta, \xi\right)-B_{n}\left(s_{\star}-\delta, \xi\right)\right) d \xi+O(\delta) .
\end{align*}
$$

We refer to the proof of [Dro19a, Lem. 2.1] for details. Hence, it suffices to estimate each summand in the RHS of (5.5).
(4) Fix $\zeta_{\star}=\left(s_{\star}, \xi_{\star}\right) \in \mathbb{Z}$. Since $H \in \mathbb{L}, \lambda_{n}(H)$ and $\lambda_{n+1}(H)$ degenerate conically at $\zeta_{\star}$. In particular, $\lambda_{n}\left(H\left(\zeta_{\star}\right)\right)>\lambda_{n-1}\left(H\left(\zeta_{\star}\right)\right)$. Therefore, $\Pi_{n-1}(\zeta)$ - hence $B(\zeta)$ depend smoothly on $\zeta$ near $\zeta_{\star}$. Using the additivity of the Berry curvature, we get

$$
\begin{equation*}
B_{n}\left(\zeta_{\star}+\varepsilon\right)=b\left(\zeta_{\star}+\varepsilon\right)+O(1) \tag{5.6}
\end{equation*}
$$

for $\varepsilon$ sufficiently small. We refer to the proof of [Dro19a, (2.21)] for details. It remains to understand $b(\zeta)$ near $\zeta_{\star}$, hence $\pi(\zeta)$ and its derivatives near $\zeta_{\star}$.
(5) Let $\left\{f_{1}, f_{2}\right\}$ be an orthonormal basis of $\operatorname{ker}\left(H\left(\zeta_{\star}\right)-\lambda_{n}\left(H\left(\zeta_{\star}\right)\right)\right.$. We define

$$
J: \mathbb{C}^{N} \longrightarrow \mathbb{C}^{2}, \quad J f=\left[\begin{array}{l}
\left\langle f, f_{1}\right\rangle \\
\left\langle f, f_{2}\right\rangle
\end{array}\right]
$$

We write a Taylor development of the $2 \times 2$ matrix $J H(\zeta) J^{*}$ near $\zeta_{\star}$ :

$$
J H\left(\zeta_{\star}+\varepsilon\right) J^{*}=J H\left(\zeta_{\star}\right) J^{*}+\sum_{j=1}^{3} B_{j} \varepsilon_{j}+O\left(\varepsilon^{2}\right)
$$

We note that $J H\left(\zeta_{\star}\right) J^{*}=\lambda_{n}\left(H\left(\zeta_{\star}\right)\right) \cdot \operatorname{Id}_{2}$ by definition of $J$. We write $B_{j}$ in the basis of Pauli matrices: $B_{j}=\sum_{k=0}^{3} a_{j k} \sigma_{k}$. This yields

$$
J H\left(\zeta_{\star}+\varepsilon\right) J^{*}=\left(\lambda_{n}\left(H\left(\zeta_{\star}\right)\right)+\sum_{j=1}^{3} a_{j 0} \varepsilon_{j}\right) \cdot \mathrm{Id}_{2}+\sum_{j, k=1}^{3} a_{j k} \sigma_{k} \varepsilon_{j}+O\left(\varepsilon^{2}\right)
$$

Let $A_{\star}$ be the $3 \times 3$ matrix with entries $a_{j k}, 1 \leqslant j, k \leqslant 3$. From Section 3.1, the eigenvalues of $J H\left(\zeta_{\star}+\varepsilon\right) J^{*}$ are

$$
\begin{equation*}
\lambda_{n}\left(H\left(\zeta_{\star}\right)\right)+\left\langle a_{0}, \varepsilon\right\rangle \pm\left|A_{\star} \varepsilon\right|+O\left(\varepsilon^{2}\right) \tag{5.7}
\end{equation*}
$$

On the other hand, the eigenvalues of $J H\left(\zeta_{\star}+\varepsilon\right) J^{*}$ are $\lambda_{n}\left(H\left(\zeta_{\star}+\varepsilon\right)\right)+O\left(\varepsilon^{2}\right)$ and $\lambda_{n+1}\left(H\left(\zeta_{\star}+\varepsilon\right)\right)+O\left(\varepsilon^{2}\right)$ - for details, see the proof of [Dro19a, (2.19)]. Since these intersect conically, $A_{\star}$ must be invertible.

For $\varepsilon \neq 0$, the matrix $\sum_{j, k=1}^{3} a_{j k} \sigma_{k} \varepsilon_{j}$ has two opposite, distinct eigenvalues. Let $\pi_{0}(\varepsilon)$ be the projector to the negative eigenvalue. Then

$$
\begin{aligned}
\pi\left(\zeta_{\star}+\varepsilon\right) & =\pi_{0}(\varepsilon)+O(|\varepsilon|), & \nabla \pi\left(\zeta_{\star}+\varepsilon\right) & =\nabla \pi_{0}(\varepsilon)+O(1) \\
\nabla \pi\left(\zeta_{\star}+\varepsilon\right) & =O\left(|\varepsilon|^{-1}\right), & \nabla \pi_{0}\left(\zeta_{\star}+\varepsilon\right) & =O\left(|\varepsilon|^{-1}\right)
\end{aligned}
$$

We refer to the proof of [Dro19a, Lem. 2.4] for such estimates. It follows that

$$
b\left(\zeta_{\star}+\varepsilon\right)=b_{0}(\varepsilon)+O\left(|\varepsilon|^{-1}\right), \quad \text { where } \quad b_{0}(\varepsilon)=\operatorname{Tr}_{\mathbb{C}^{N}}\left(\pi_{0}(\varepsilon)\left[\partial_{\xi_{1}} \pi_{0}(\varepsilon), \partial_{\xi_{2}} \pi_{0}(\varepsilon)\right]\right)
$$

Grouping with (5.6), we obtain $B_{n}\left(\zeta_{\star}+\varepsilon\right)=b_{0}(\varepsilon)+O\left(|\varepsilon|^{-1}\right)$. In particular,

$$
\begin{align*}
\int_{\left|\xi-\xi_{\star}\right| \leqslant r} B_{n}\left(s_{\star} \pm \delta, \xi\right) d \xi & =\int_{\left|\xi-\xi_{\star}\right| \leqslant r} b_{0}\left( \pm \delta, \xi-\xi_{\star}\right) d \xi+O\left(\int_{\left|\xi-\xi_{\star}\right| \leqslant r} \frac{1}{\left|\xi-\xi_{\star}\right|} d \xi\right)  \tag{5.8}\\
& =\int_{|\xi| \leqslant r} b_{0}( \pm \delta, \xi) d \xi+O(r)
\end{align*}
$$

(6) Since $A_{\star}$ is invertible, the three matrices $A_{j}=\sum_{k=1}^{3} a_{j k} \sigma_{k}, 1 \leqslant j \leqslant 3$, form a basis of $\mathcal{E}_{0}$. We apply the Gran-Schmidt process to $\left(A_{1}, A_{2}, A_{3}\right)$ : there exists $\left(\widetilde{\sigma}_{1}, \widetilde{\sigma}_{2}, \widetilde{\sigma}_{3}\right)$ orthonormal basis of $\mathcal{E}_{0}$ and $\left(t_{j k}\right) \in M_{3}(\mathbb{R})$ upper triangular with positive elements on the diagonal such that $A_{j}=\sum_{k=1}^{3} t_{j k} \widetilde{\sigma}_{k}$.

We write $\widetilde{\sigma}_{k}=\epsilon_{\star} \cdot U \sigma_{k} U^{*}$, where $\epsilon_{\star}$ is the determinant of ( $\left.\widetilde{\sigma}_{1}, \widetilde{\sigma}_{2}, \widetilde{\sigma}_{3}\right)$ with respect to $\left(\sigma_{1}, \sigma_{2}, \sigma_{3}\right)-\operatorname{see}(5.3)$. In particular, $\epsilon_{\star}=\operatorname{sgn}\left(\operatorname{det}\left(A_{\star}\right)\right)$. It follows that

$$
A_{j}=\epsilon_{\star} \cdot U\left(\sum_{k=1}^{3} t_{j k} \sigma_{k}\right) U^{*}, \quad \sum_{j, k=1}^{3} a_{j k} \sigma_{k} \varepsilon_{j}=\epsilon_{\star} \cdot U\left(\sum_{j, k=1}^{3} t_{j k} \sigma_{k} \varepsilon_{j}\right) U^{*}
$$

Hence, $\pi_{0}(\varepsilon)$ is, up to conjugation, the projector associated to the negative eigenvalue of $\epsilon_{\star} \cdot \sum_{j, k=1}^{3} t_{j k} \sigma_{k} \varepsilon_{j}$.

We define more appropriate coordinates

$$
\begin{equation*}
\widetilde{\xi}_{1}=\frac{t_{12} \delta+t_{22} \xi_{1}}{t_{11} \delta}, \quad \widetilde{\xi}_{2}=\frac{t_{13} \delta+t_{23} \xi_{1}+t_{33} \xi_{2}}{t_{11} \delta} \tag{5.9}
\end{equation*}
$$

$\underset{\sim}{\text { Using }} \underset{\sim}{\tilde{\xi}}) d$ ariance of two-forms under change of coordinates, $b_{0}(\xi) d \xi=\widetilde{b}_{0}(\widetilde{\xi}) d \widetilde{\xi}$, where $\widetilde{b}_{0}( \pm \delta, \widetilde{\xi}) d \xi$ is the two-form associated to the negative eigenspace of

$$
\epsilon_{\star} \delta \cdot t_{11}\left(\sigma_{1}+\sigma_{2} \widetilde{\xi}_{1}+\sigma_{3} \widetilde{\xi}_{2}\right)
$$

This setup allows us to apply [FC13, (23)], which gives:

$$
\widetilde{b}_{0}( \pm \delta, \widetilde{\xi})=\frac{i \epsilon_{\star}^{3}( \pm \delta)^{3}}{2 \delta^{3}\left(\widetilde{\xi}_{1}^{2}+\widetilde{\xi}_{2}^{2}+1\right)^{3 / 2}}\left[\begin{array}{c}
1 \\
\widetilde{\xi}_{1} \\
\widetilde{\xi}_{2}
\end{array}\right] \cdot\left[\begin{array}{l}
0 \\
1 \\
0
\end{array}\right] \wedge\left[\begin{array}{l}
0 \\
0 \\
1
\end{array}\right]=\frac{ \pm \epsilon_{\star}}{2\left(\widetilde{\xi}_{1}^{2}+\widetilde{\xi}_{2}^{2}+1\right)^{3 / 2}}
$$

Under the change of coordinates (5.9), the disk $|\xi| \leqslant r$ gets mapped to an ellipse centered at distance $O(1)$ from the origin, of dimensions $\sim \delta^{-1}$. Thus,

$$
\begin{equation*}
\int_{|\xi| \leqslant r} b_{0}( \pm \delta, \xi) d \xi=\int_{\mathbb{R}^{2}} b_{0}( \pm \delta, \xi) d \xi+O(\delta)= \pm \epsilon_{\star} \pi+O(\delta) \tag{5.10}
\end{equation*}
$$

We refer to the proof of [Dro19a, Lem. 2.5] for details.
(7) Grouping (5.4), (5.5), (5.8) and (5.10), we end up with

$$
c_{1}\left(\mathscr{V}_{1}\right)-c_{1}\left(\mathscr{V}_{0}\right)=\sum_{\zeta_{\star} \in \mathcal{R}} \epsilon_{\star}+O(r+\delta)=\sum_{\zeta_{\star} \in \mathcal{R}} \operatorname{sgn}\left(\operatorname{det}\left(A_{\star}\right)\right)+O(r+\delta)
$$

Making $\delta \rightarrow 0$, we end up with

$$
\begin{equation*}
c_{1}\left(\mathscr{V}_{1}\right)-c_{1}\left(\mathscr{V}_{0}\right)=\sum_{\zeta_{\star} \in \mathcal{R}} \operatorname{sgn}\left(\operatorname{det}\left(A_{\star}\right)\right)+O(r) \tag{5.11}
\end{equation*}
$$

Taking $r$ sufficiently small, the term $O(r)$ is at most $1 / 2$. Since both sides of (5.11) are integers, we end up with

$$
c_{1}\left(\mathscr{V}_{1}\right)-c_{1}\left(\mathscr{V}_{0}\right)=\sum_{\zeta_{\star} \in \mathcal{R}} \operatorname{sgn}\left(\operatorname{det}\left(A_{\star}\right)\right)
$$

This completes the proof.

## Appendix. Continuous approximation

Let $H \in \mathcal{L}$ with a conical degeneracy at $\left(s_{0}, \xi_{0}\right) \in(0,1) \times \mathbb{T}^{2}$ and $\mathcal{H}^{\delta}$ defined as in (1.7). In this appendix, we derive formally the effective Dirac equation (1.10). It describes the evolution of amplitudes to solutions of $\left(D_{t}-\mathcal{H}^{\delta}\right) \psi=0$ that are initially concentrated (in phase-space) near $\left(\mathbb{R} e_{1}+s_{0} e_{2}, \xi_{0}\right)$.
A.1. Reduction to $\left(s_{0}, \xi_{0}\right)=(0,0)$. - We show that $\mathcal{H}^{\delta}$ is unitarily equivalent to an operator with a conical degeneracy at $(0,0)$. Define

$$
\widetilde{H}_{s}(\xi)=H_{s_{0}+s}\left(\xi+\xi_{0}\right), \quad \widetilde{H}_{s-s_{0}}\left(\xi-\xi_{0}\right)=H_{s}(\xi)
$$

and $\widetilde{\mathcal{H}}_{s}, \widetilde{\mathrm{Q}}^{\delta}, \widetilde{\mathcal{H}}^{\delta}$ relative to $\widetilde{H}$, according to (1.4) and (1.7).
For $m \in \mathbb{Z}^{2}$, set $\ell=m-\left[\delta^{-1} s_{0}\right] e_{2}$, where $\left[\delta^{-1} s_{0}\right]$ stands for the integer part of $\delta^{-1} s_{0}$. For $\phi \in \ell^{2}\left(\mathbb{Z}^{2}, \mathbb{C}^{N}\right)$, we have:

$$
\begin{aligned}
e^{i \xi_{0} \ell} \cdot\left(\widetilde{\mathbb{Q}}^{\delta} \phi\right)(\ell) & =e^{i \xi_{0} \ell} \cdot\left(\widetilde{Q}_{\delta n_{2}-s_{0}} \phi\right)(\ell) \\
& =\int_{\mathbb{T}^{2}} e^{i\left(\xi+\xi_{0}\right) \ell} \cdot \widetilde{H}_{\delta n_{2}-s_{0}}(\xi) \widehat{\phi}(\xi) \cdot \frac{d \xi}{(2 \pi)^{2}} \\
& =\int_{\mathbb{T}^{2}} e^{i \xi \ell} \cdot \widetilde{H}_{\delta n_{2}-s_{0}}\left(\xi-\xi_{0}\right) \widehat{\phi}\left(\xi-\xi_{0}\right) \cdot \frac{d \xi}{(2 \pi)^{2}} \\
& =\int_{\mathbb{T}^{2}} e^{i \xi \ell} \cdot H_{\delta n_{2}}(\xi) \widehat{e^{i \xi_{0} \cdot \phi}}(\xi) \cdot \frac{d \xi}{(2 \pi)^{2}}=\left(\mathbb{Q}^{\delta} e^{i \xi_{0} \cdot} \phi\right)(\ell) .
\end{aligned}
$$

This means that $U \widetilde{\mathfrak{Q}}^{\delta} U^{*}=\mathbb{Q}^{\delta}$, where

$$
U \phi(m)=\left(e^{i \xi_{0} \cdot} \phi\right)\left(m-\left[\delta^{-1} s_{0}\right] e_{2}\right), \quad U^{*} \phi(m)=e^{i \xi_{0} m} \cdot \phi\left(m+\left[\delta^{-1} s_{0}\right] e_{2}\right)
$$

This implies $U \widetilde{\mathcal{H}}^{\delta} U^{*}=\mathcal{H}^{\delta}: \widetilde{\mathcal{H}}^{\delta}$ and $\mathcal{H}^{\delta}$ are unitarily equivalent.
A.2. Effective equation. - Since $H_{s}(\xi)$ has a conical degeneracy at $\left(s_{0}, \xi_{0}\right)$, there exists $f_{1}, f_{2} \in \mathbb{C}^{N}$ satisfying (1.5). As $\delta \rightarrow 0$, we derive (formally) the leading asymptotics of $\mathcal{H}^{\delta} \phi=\mathbb{Q}^{\delta} \phi+O(\delta)$, where
(A.1) $\quad \phi(m)=e^{i \xi_{0} m} \cdot \sum_{j=1}^{2} \alpha_{j}\left(s_{0} e_{2}+\delta^{1 / 2} m\right) f_{j} \in \ell^{2}\left(\mathbb{Z}^{2}, \mathbb{C}^{N}\right), \quad \alpha \in C_{0}^{\infty}\left(\mathbb{R}^{2}, \mathbb{C}^{N}\right)$.

After rescaling, $\phi$ is semiclassically (scale $\delta$ ) localized near $\left(\mathbb{R} e_{1}+s_{0} e_{2}, \xi_{0}\right)$. We write (A.1) as $\phi \simeq U^{*} \varphi$, where $\varphi(m)=\sum_{j=1}^{2} \alpha_{j}\left(\delta^{1 / 2} m\right) f_{j}=J^{*} \alpha\left(\delta^{1 / 2} m\right)$, and $J: \mathbb{C}^{N} \rightarrow \mathbb{C}^{2}$ is the operator of (1.5). Using a Riemann sum argument, we observe that as $\delta \rightarrow 0$,

$$
\begin{align*}
\delta \cdot \widehat{\varphi}\left(\delta^{1 / 2} \xi\right) & =J^{*}\left(\delta \sum_{m \in \mathbb{Z}^{2}} e^{-i \delta^{1 / 2} \xi m} \alpha\left(\delta^{1 / 2} m\right)\right)  \tag{A.2}\\
& \simeq J^{*}\left(\int_{\mathbb{R}^{2}} e^{-i \xi x} \alpha(x) d x\right)=J^{*} \widehat{\alpha}(\xi)
\end{align*}
$$

Thanks to Section A.1, we have $\mathcal{H}^{\delta} \phi=\mathcal{H}^{\delta} U^{*} \varphi=U^{*} \widetilde{\mathcal{H}}^{\delta} \varphi$. Now, we compute $\widetilde{\mathcal{H}}^{\delta} \varphi$ :

$$
\begin{align*}
\left(\widetilde{\mathcal{H}}^{\delta} \varphi\right)(m) & \simeq\left(\widetilde{\mathbb{Q}}^{\delta} \varphi\right)(m)=\left(\widetilde{\mathcal{H}}_{\delta m_{2}} \varphi\right)(m)=\int_{\mathbb{T}^{2}} e^{i \xi m} \cdot \widetilde{H}_{\delta m_{2}}(\xi) \widehat{\varphi}(\xi) \cdot \frac{d \xi}{(2 \pi)^{2}} \\
& =\int_{\delta^{-1} \mathbb{T}^{2}} e^{i \delta^{1 / 2} \xi m} \cdot H_{\delta m_{2}}\left(\delta^{1 / 2} \xi\right) \widehat{\varphi}\left(\delta^{1 / 2} \xi\right) \cdot \frac{\delta d \xi}{(2 \pi)^{2}}, \tag{A.3}
\end{align*}
$$

where we made the substitution $\xi \mapsto \delta^{1 / 2} \xi$. Since $\varphi$ is spectrally concentrated near 0 , it is reasonable to replace the integration domain in (A.3) to $\mathbb{R}^{2}$. Using (A.2), we get

$$
\begin{equation*}
\left(\widetilde{\mathcal{H}}^{\delta} \varphi\right)(m) \simeq \int_{\mathbb{R}^{2}} e^{i \delta^{1 / 2} \xi m} \cdot H_{\delta m_{2}}\left(\delta^{1 / 2} \xi\right) J^{*} \widehat{\alpha}(\xi) \cdot \frac{d \xi}{(2 \pi)^{2}} \tag{A.4}
\end{equation*}
$$

The identity (1.5) allows us to expand $H_{\delta m_{2}}\left(\delta^{1 / 2} \xi\right) J^{*}$ as

$$
\begin{equation*}
H_{\delta m_{2}}\left(\delta^{1 / 2} \xi\right) J^{*} \simeq J^{*}\left(E_{0}+\delta^{1 / 2} \cdot \not D\left(\delta^{1 / 2} m_{2}, \xi\right)\right) \tag{A.5}
\end{equation*}
$$

where $\not D(s, \xi)$ is a family of $2 \times 2$ matrices depending linearly on $(s, \xi)$, and $E_{0}$ is the energy of the conical crossing. Plugging (A.5) into (A.4), we obtain

$$
\begin{aligned}
\left(\widetilde{\mathcal{H}}^{\delta} \varphi\right)(m) & \simeq J^{*} \int_{\mathbb{R}^{2}} e^{i \delta^{1 / 2} \xi m} \cdot\left(E_{0}+\delta^{1 / 2} \not D\left(\delta^{1 / 2} m_{2}, \xi\right)\right) \widehat{\alpha}(\xi) \cdot \frac{d \xi}{(2 \pi)^{2}} \\
& =J^{*}\left(E_{0}+\delta^{1 / 2} \not D\right) \alpha\left(\delta^{1 / 2} m\right)
\end{aligned}
$$

where $\not D=\not D\left(x_{2}, D_{x}\right)$ is a Dirac operator. Since $\varphi(m)=J^{*} \alpha\left(\delta^{1 / 2} m\right)$, this means that $J^{*}$ approximately intertwines between $\widetilde{\mathcal{H}}^{\delta}$ and $E_{0}+\delta^{1 / 2} \not D$, for adiabatic data.

Up to a phase and a time-rescaling, the equations $D_{t}-E_{0}-\delta^{1 / 2} \not D$ and $D_{t}-\not D$ are equivalent. Using the above intertwining, we conclude that $\left(D_{t}-\mathcal{H}^{\delta}\right) \psi=0$ has approximate solutions whose asymptotics are slow linear combinations of $f_{1}$ and $f_{2}$ :

$$
e^{i\left(E_{0} t+\xi_{0} m\right)} \cdot \sum_{j=1}^{2} \beta_{j}\left(\delta^{1 / 2} t, s_{0} e_{2}+\delta^{1 / 2} m\right) f_{j}
$$

with amplitudes $\beta_{j}(t, x)$ solving the Dirac equation (1.10): $\left(D_{t}-\not D\right) \beta=0$.

## References

[ADH20] H. Ammari, B. Davies \& E. O. Hiltunen - "Robust edge modes in dislocated systems of subwavelength resonators", 2020, arXiv:2001.10455.
[ADHY19] H. Ammari, B. Davies, E. O. Hiltunen \& S. Yu - "Topologically protected edge modes in one-dimensional chains of subwavelength resonators", 2019, arXiv:1906.10688.
$\left[\mathrm{AFH}^{+} 18\right]$ H. Аmmari, B. Fitzpatrick, E. O. Hiltunen, H. Lee \& S. Yu - "Honeycomb-lattice Minnaert bubbles", 2018, arXiv:1811.03905.
[Arn95] V. I. Arnold - "Remarks on eigenvalues and eigenvectors of Hermitian matrices, Berry phase, adiabatic connections and quantum Hall effect", Selecta Math. (N.S.) 1 (1995), no. 1, p. 1-19.
[ASBVB13] J. C. Avila, H. Schulz-Baldes \& C. Villegas-Blas - "Topological invariants of edge states for periodic two-dimensional models", Math. Phys. Anal. Geom. 16 (2013), no. 2, p. 137170.
[AS78] J. E. Avron \& B. Simon - "Analytic properties of band functions", Ann. Physics 110 (1978), no. 1, p. 85-101.
[Bal19a] G. BaL - "Continuous bulk and interface description of topological insulators", J. Math. Phys. 60 (2019), no. 8, p. 081506, 20.
[Bal19b] , "Topological protection of perturbed edge states", Commun. Math. Sci. 17 (2019), no. 1, p. 193-225.
[Bal19c] , "Topological invariants for interface modes", 2019, arXiv:1906.08345.
[Ber84] M. V. Berry - "Quantal phase factors accompanying adiabatic changes", Proc. Roy. Soc. London Ser. A 392 (1984), no. 1802, p. 45-57.
[BKR17] C. Bourne, J. Kellendonk \& A. Rennie - "The $K$-theoretic bulk-edge correspondence for topological insulators", Ann. Inst. H. Poincaré Phys. Théor. 18 (2017), no. 5, p. 18331866.
[Col91] Y. Colin de Verdière - "Sur les singularités de van Hove génériques", in Analyse globale et physique mathématique (Lyon, 1989), Mém. Soc. Math. France (N.S.), vol. 46, Société Mathématique de France, Paris, 1991, p. 99-110.
[DE99] L. Dieci \& T. Eirola - "On smooth decompositions of matrices", SIAM J. Matrix Anal. Appl. 20 (1999), no. 3, p. 800-819.
[DP12] L. Dieci \& A. Pugliese - "Hermitian matrices depending on three parameters: coalescing eigenvalues", Linear Algebra Appl. 436 (2012), no. 11, p. 4120-4142.
[Dom11] М. Dомокоs - "Discriminant of symmetric matrices as a sum of squares and the orthogonal group", Comm. Pure Appl. Math. 64 (2011), no. 4, p. 443-465.
[Dro19a] A. Drouot - "The bulk-edge correspondence for continuous honeycomb lattices", Comm. Partial Differential Equations 44 (2019), no. 12, p. 1406-1430.
[Dro19b] , "Characterization of edge states in perturbed honeycomb structures", Pure Appl. Anal. 1 (2019), no. 3, p. 385-445.
[Dro19c] , "Microlocal analysis of the bulk-edge correspondence", arXiv:1909.10474, 2019.
[DW20] A. Drouot \& M. I. Weinstein - "Edge states and the valley Hall effect", Adv. Math. 368 (2020), p. 107142, 51.
[EGS05] A. Elgart, G. M. Graf \& J. H. Schenker - "Equality of the bulk and edge Hall conductances in a mobility gap", Comm. Math. Phys. 259 (2005), no. 1, p. 185-221.
[FLTW16] C. L. Fefferman, J. P. Lee-Thorp \& M. I. Weinstein - "Edge states in honeycomb structures", Ann. PDE 2 (2016), no. 2, article no. 12 (80 pages).
[FLTW18] C. L. Fefferman, J. P. Lee-Thorp \& M. I. Weinstein - "Honeycomb Schrödinger operators in the strong binding regime", Comm. Pure Appl. Math. 71 (2018), no. 6, p. 1178-1270.
[FW12] C. L. Fefferman \& M. I. Weinstein - "Honeycomb lattice potentials and Dirac points", J. Amer. Math. Soc. 25 (2012), no. 4, p. 1169-1220.
[FK04] C. Fermanian Kammerer - "Semiclassical analysis of generic codimension 3 crossings", Internat. Math. Res. Notices (2004), no. 45, p. 2391-2435.
[FKG03] C. Fermanian Kammerer \& P. Gérard - "A Landau-Zener formula for non-degenerated involutive codimension 3 crossings", Ann. Inst. H. Poincaré Phys. Théor. 4 (2003), no. 3, p. 513-552.
[FT16] S. Freund \& S. Teufel - "Peierls substitution for magnetic Bloch bands", Ann. PDE 9 (2016), no. 4, p. 773-811.
[FC13] M. Fruchart \& D. Carpentier - "An introduction to topological insulators", Comptes Rendus Physique 14 (2013), no. 9, p. 779-815.
[GP13] G. M. Graf \& M. Porta - "Bulk-edge correspondence for two-dimensional topological insulators", Comm. Math. Phys. 324 (2013), no. 3, p. 851-895.
[GP74] V. Guillemin \& A. Pollack - Differential topology, Prentice-Hall, Inc., Englewood Cliffs, N.J., 1974.
[Hat93] Y. Hatsugai - "Chern number and edge states in the integer quantum Hall effect", Phys. Rev. Lett. 71 (1993), no. 22, p. 3697-3700.
[Ily92] N. V. Llyushechkin - "The discriminant of the characteristic polynomial of a normal matrix", Mat. Zametki 51 (1992), no. 3, p. 16-23, 143.
[KRSB02] J. Kellendonk, T. Richter \& H. Schulz-Baldes - "Edge current channels and Chern numbers in the integer quantum Hall effect", Rev. Math. Phys. 14 (2002), no. 1, p. 87119.
[Kuc16] P. Kuchment - "An overview of periodic elliptic operators", Bull. Amer. Math. Soc. (N.S.) 53 (2016), no. 3, p. 343-414.
[KP07] P. Kuchment \& O. Post - "On the spectra of carbon nano-structures", Comm. Math. Phys. 275 (2007), no. 3, p. 805-826.
[Lax98] P. D. Lax - "On the discriminant of real symmetric matrices", Comm. Pure Appl. Math 51 (1998), no. 11-12, p. 1387-1396.
[Lee16] M. Lee - "Dirac cones for point scatterers on a honeycomb lattice", SIAM J. Math. Anal. 48 (2016), no. 2, p. 1459-1488.
[LTWZ19] J. P. Lee-Thorp, M. I. Weinstein \& Y. Zhu - "Elliptic operators with honeycomb symmetry: Dirac points, edge states and applications to photonic graphene", Arch. Rational Mech. Anal. 232 (2019), no. 1, p. 1-63.
[Mac15] I. G. Macdonald - Symmetric functions and Hall polynomials, 2nd ed., Oxford Classic Texts in the Physical Sciences, The Clarendon Press, Oxford University Press, New York, 2015.
[Mon17] D. Monaco - "Chern and Fu-Kane-Mele invariants as topological obstructions", in Advances in quantum mechanics, Springer INdAM Ser., vol. 18, Springer, Cham, 2017, p. 201-222.
[MP14] D. Monaco \& G. Panati - "Topological invariants of eigenvalue intersections and decrease of Wannier functions in graphene", J. Statist. Phys. 155 (2014), no. 6, p. 1027-1071.
[Moo01] J. D. Moore - Lectures on Seiberg-Witten invariants, 2nd ed., Lect. Notes in Math., vol 1629, Springer-Verlag, Berlin, 2001.
[vNW29] J. von Neumann \& E. Wigner - "Über das Verhalten von Eigenwerten bei adiabatischen Prozessen", Phys. Z. 30 (1929), p. 467-470.
[Pan07] G. Panati - "Triviality of Bloch and Bloch-Dirac bundles", Ann. Inst. H. Poincaré Phys Théor. 8 (2007), no. 5, p. 995-1011.
[PST03] G. Panati, H. Spohn \& S. Teufel - "Effective dynamics for Bloch electrons: Peierls substitution and beyond", Comm. Math. Phys. 242 (2003), no. 3, p. 547-578.
[Par02] B. N. Parlett - "The (matrix) discriminant as a determinant", Linear Algebra Appl. $\mathbf{3 5 5}$ (2002), p. 85-101.
[Pet16] P. Petersen - Riemannian geometry, 3rd ed., Graduate Texts in Math., vol. 171, Springer, Cham, 2016.
[PSB16] E. Prodan \& H. Schulz-Baldes - Bulk and boundary invariants for complex topological insulators, Mathematical Physics Studies, Springer, Cham, 2016.
[RH08] S. Raghu \& F. D. M. Haldane - "Analogs of quantum-Hall-effect edge states in photonic crystals", Phys. Rev. A $\mathbf{7 8}$ (2008), article no. 033834 (21 pages).
[See64] R. T. Seeley - "Extension of $C^{\infty}$ functions defined in a half space", Proc. Amer. Math. Soc. 15 (1964), p. 625-626.
[Ser10] D. Serre - Matrices. Theory and applications, 2nd ed., Graduate Texts in Math., vol 216, Springer, New York, 2010.
[Sim83] B. Simon - "Holonomy, the quantum adiabatic theorem, and Berry's phase", Phys. Rev Lett. 51 (1983), no. 24, p. 2167-2170.
[Sin05] S. F. Singer - Linearity, symmetry, and prediction in the hydrogen atom, Undergraduate Texts in Math., Springer, New York, 2005.
[Tey99] M. Teytel - "How rare are multiple eigenvalues?", Comm. Pure Appl. Math. 52 (1999), no. 8, p. 917-934.
[Wal47] P. R. Wallace - "The band theory of graphite", Phys. Rev., II. Ser. 71 (1947), p. 622634.

Manuscript received Ifth May 2020
accepted 8th February 2021
Alexis Drouot, Padelford Hall, University of Washington
W Stevens Way NE, Seattle, WA 98105, USA
E-mail : alexis.drouot@gmail.com


[^0]:    Mathematical subject classification (2020). - 47A13, 81Q10, 81Q05.
    Keywords. - Dirac cones, topological insulators, Chern numbers.

    I thankfully acknowledge support from NSF DMS-1440140 (MSRI, Fall 2019) and DMS-1800086, and from the Simons Foundation through M. I. Weinstein's Math+X investigator award \#376319.

