# Bifurcation Analysis in a Frustrated Nematic Cell 

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#### Abstract

Using Landau-de Gennes theory to describe nematic order, we study a frustrated cell consisting of nematic liquid crystal confined between two parallel plates. We prove the uniqueness of equilibrium states for a small cell width. Letting the cell width grow, we study the behavior of this unique solution. Restricting ourselves to a certain interval of temperature, we prove that this solution becomes unstable at a critical value of the cell width. Moreover, we show that this loss of stability comes with the appearance of two new solutions: there is a symmetric pitchfork bifurcation. This picture agrees with numerical simulations performed by Palffy-Muhoray, Gartland and Kelly, and also by Bisi, Gartland, Rosso, and Virga. Some of the methods that we use in the present paper apply to other situations, and we present the proofs in a general setting. More precisely, the paper contains the proof of a general uniqueness result for a class of perturbed quasilinear elliptic systems, and general considerations about symmetric solutions and their stability, in the spirit of Palais' Principle of Symmetric Criticality.


Keywords Bifurcation • Liquid crystals • Hybrid cell
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## 1 Introduction

In a nematic liquid crystal, rigid rod-like molecules tend to align in a common preferred direction. To describe this orientational order, de Gennes (De Gennes and Prost 1993)

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Fig. 1 Schematic representation of EE and BD configurations: variations of the eigenframe of $Q$ through the cell. Eigenvectors corresponding to the largest eigenvalues are emphasized. In an EE configuration, as we move through the cell from the left to the right, the eigenvalue associated to $\mathbf{e}_{\mathbf{z}}$ (the prescribed director on the left plate) decreases until, in the middle of the cell, it becomes equal to the eigenvalue associated to $\mathbf{e}_{\mathbf{y}}$ (the prescribed director on the right plate), which in turn increases to match the boundary condition. In a BD configuration, as we move through the cell from the left to the right, the eigenframe rotates, so that the eigenvector corresponding to the largest eigenvalue rotates from $\mathbf{e}_{\mathbf{z}}$ to $\mathbf{e}_{\mathbf{y}}$. Note that a similar eigenframe rotation could occur in the opposite way as the one pictured here: there are two possible types of BD configurations
introduced the so-called $Q$-tensor: a $3 \times 3$ traceless symmetric matrix. The eigenframe of the $Q$-tensor describes the principal mean directions of alignment, while the corresponding eigenvalues describe the degrees of alignment along those directions. A null $Q$-tensor corresponds to the isotropic liquid state. A $Q$-tensor with two equal eigenvalues corresponds to the uniaxial state, which is axially symmetric around one eigenvector, called the director. The generic case of a $Q$-tensor with three distinct eigenvalues corresponds to the biaxial state. There exist other continuum theories for the description of nematic order. The Oseen-Frank theory and the Ericksen theory make use of simpler order parameters but do not account for biaxiality. The $Q$-tensor description allows for a much finer understanding of some phenomena, as defects (see e.g., Sonnet et al. 1995) and material frustration (Palffy-Muhoray 1994).

In the present paper we focus on a hybrid cell consisting of nematic material confined between two parallel bounding plates, and subject to competing strong anchoring conditions on each plate. Hybrid-aligned nematic cells are interesting for technological applications (Huang et al. 2011; Matsumoto et al. 1976), as is, in general, the effect of confining geometries on liquid crystals (Chiccoli et al. 2003). Nematic systems similar to the one considered here have been studied numerically in Bisi et al. (2003) and Palffy-Muhoray (1994) as a model for material frustration. On each bounding plate, the prescribed boundary condition is uniaxial and parallel to the plate, with director orthogonal to the one prescribed on the opposite plate. The numerics presented in Bisi et al. (2003) and Palffy-Muhoray (1994) bring to light two different families of solutions: eigenvalue exchange (EE) and bent director (BD) configurations. In an eigenvalue exchange solution, the $Q$-tensor eigenframe remains constant through the whole cell, and only the eigenvalues vary to match the boundary conditions. Therefore, inside the cell the material is strongly biaxial. In the BD configuration, however, the eigenframe rotates to connect the two orthogonal uniaxial states on the plates. Hence the tensor remains approximately uniaxial, with director bending from one plate to the other. Those two kinds of configurations are depicted in Fig. 1.

Fig. 2 Shape of the pitchfork bifurcation described in Bisi et al. (2003) and Palffy-Muhoray (1994)


When working with dimensionless variables, two parameters influence the behavior of the system: a reduced temperature $\theta$, and a typical length $\lambda$ proportional to the thickness of the cell. In Bisi et al. (2003) and Palffy-Muhoray (1994), a bifurcation analysis is performed numerically as the cell thickness varies, at fixed temperatures. In both studies, a symmetric pitchfork bifurcation diagram is obtained (Bisi et al. 2003, Fig. 8), which can be described as follows (see Fig. 2). When the parameter $\lambda$ (and thus the cell thickness) is small, the only equilibrium is an eigenvalue exchange configuration, which is stable. Letting the cell thickness grow, a critical value is attained, at which this eigenvalue exchange solution loses stability. At this point, bifurcation occurs and two new stable branches of solutions appear, corresponding to BD configurations, with their eigenframe rotating in one way or the other. The results pictured in Bisi et al. (2003, Fig. 8) were obtained for a special value of the reduced temperature, $\theta=-8$, at which computations are simplified.

In the present paper, we aim at providing rigorous mathematical arguments justifying the shape of the bifurcation diagram pictured in Fig. 2. Thus, we fix the temperature and let the cell thickness vary.

In a first step, we study the limits of small and large cell thickness. For a very large cell thickness, we check that energy minimizers converge toward two possible limiting uniaxial configurations, corresponding to a rotation of the director in one way or the other. On the other hand, when the cell is sufficiently narrow, we prove indeed that the energy admits a unique critical point. Symmetry considerations imply that this unique solution is an eigenvalue exchange configuration-thus showing that Fig. 2 is valid for small $\lambda$. The method used to prove uniqueness applies to a quite wide class of problems, and we prove a general uniqueness result for a class of perturbed quasilinear elliptic systems in Appendix 2.

In a second step, we perform a bifurcation analysis and show that there is indeed a symmetric pitchfork bifurcation, at least when the reduced temperature $\theta$ is close to $\theta=-8$ (the special value at which Bisi et al. (2003, Fig. 8) was obtained). More specifically, we prove the following result.

Theorem Let $\theta \approx-8$. Consider, for small $\lambda$, the unique solution $\chi_{\lambda}$. The branch of eigenvalue exchange solutions $\lambda \mapsto \chi_{\lambda}$ may be extended smoothly to larger $\lambda$, and loses stability at a critical value $\lambda_{c}$. At this point, a symmetric pitchfork bifurcation occurs.

More precisely, we prove first the above Theorem in the case $\theta=-8$. Then, we identify the properties that make this special case work, which leads to an abstract result
of the form: if $\theta$ satisfies some properties, then bifurcation occurs. And eventually we check that those properties are stable: if a $\theta_{0}$ satisfies them, they extend to nearby $\theta \approx \theta_{0}$. In particular we obtain the above Theorem.

Similar nematic systems have also been studied using Oseen-Frank theory (Barbero and Barberi 1983) and Ericksen theory (Ambrosio and Virga 1991). Oseen-Frank theory describes nematic order with help of a sole director field, while in Ericksen theory the degree of alignment along the director is also accounted for. In both cases, only uniaxial states can be described, so that there is no equivalent to the eigenvalue exchange solution considered in the present paper. However, the aforementioned works (Ambrosio and Virga 1991; Barbero and Barberi 1983) do bring into light the existence of a critical cell width-as in our case-below which the BD solution is no longer valid. Within Oseen-Frank theory (Barbero and Barberi 1983), the boundary anchoring is relaxed and a uniform configuration is preferred below the critical width. Within Ericksen theory (Ambrosio and Virga 1991), the strong boundary anchoring is preserved, forcing the director field to have a discontinuity in the cell center. In our case such a discontinuity is avoided by the eigenvalue exchange mechanism.

The plan of the paper is the following. In Sect. 2, we present the precise model used to describe the cell. In Sect. 3, we discuss the existence and some properties of eigenvalue exchange configurations. In Sect. 4, we study the limits of large and small cell thickness. Then, we concentrate on the unique branch of eigenvalue exchange solutions starting from small $\lambda$, and show that a symmetric pitchfork bifurcation occurs. We treat the case $\theta=-8$ in Sect. 5, and the perturbed case $\theta \approx-8$ in Sect. 6 .

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## 2 Model

The cell consists of nematic material confined between two parallel bounding plates (of infinite size), with competing strong anchoring conditions on each plate. In an orthonormal basis $\left(\overrightarrow{e_{x}}, \overrightarrow{e_{y}}, \overrightarrow{e_{z}}\right)$, the bounding plates are perpendicular to $\overrightarrow{e_{x}}$ and parallel to the $(y, z)$ plane. The width of the cell is $2 d$ : one plate at $x=-d$, the other at $x=d$. On the left plate $(x=-d)$, the boundary condition is uniaxial with director $\overrightarrow{e_{z}}$, and on the right plate $(x=d)$ the boundary condition is uniaxial with director $\overrightarrow{e_{y}}$ (see Fig. 1).

Nematic order is described by means of de Gennes' $Q$-tensor-a traceless symmetric $3 \times 3$ matrix -, and Landau-de Gennes free energy density

$$
e(Q)=\frac{L}{2}|\nabla Q|^{2}+f_{b}(Q),
$$

where the bulk energy density $f_{b}$ is given by

$$
f_{b}(Q)=\frac{a(T)}{2}|Q|^{2}-\frac{b}{3} \operatorname{tr}\left(Q^{3}\right)+\frac{c}{4}|Q|^{4} .
$$

Here and in the sequel, the notation $|\cdot|$ refers to the usual euclidean norm $|Q|^{2}=$ $\operatorname{tr}\left(Q^{t} Q\right)=\sum Q_{i j} Q_{i j}$.

In Bisi et al. (2003) and Palffy-Muhoray (1994), the numerical simulations are performed under two symmetry restrictions : the $Q$-tensor depends only on $x$, and $\overrightarrow{e_{x}}$ is always an eigenvector. These restrictions are natural, since the system is invariant in the $x$ and $y$ directions, and since $\overrightarrow{e_{x}}$ is an eigenvector of the boundary conditions. It is not our goal here to justify rigorously the validity of these symmetry assumptions : we will, from the beginning, consider $Q$-tensors depending only on $x$, with $\overrightarrow{e_{x}}$ as an eigenvector. The assumption on the eigenvector implies that the $Q$-tensors considered here will always have four zero entries, as in (1) below.

More precisely, we will study maps

$$
Q(x)=\left(\begin{array}{ccc}
-2 q_{1}(x) & 0 & 0  \tag{1}\\
0 & q_{1}(x)-q_{2}(x) & q_{3}(x) \\
0 & q_{3}(x) & q_{1}(x)+q_{2}(x)
\end{array}\right), \quad x \in[-d, d]
$$

minimizing the energy functional

$$
E(Q)=\int_{-d}^{d}\left(\frac{L}{2}\left|Q^{\prime}\right|^{2}+f_{b}(Q)\right) \mathrm{d} x
$$

when subject to boundary conditions

$$
Q(-d)=\left(\begin{array}{ccc}
-2 q_{+} & 0 & 0 \\
0 & -2 q_{+} & 0 \\
0 & 0 & 4 q_{+}
\end{array}\right), \quad Q(-d)=\left(\begin{array}{ccc}
-2 q_{+} & 0 & 0 \\
0 & 4 q_{+} & 0 \\
0 & 0 & -2 q_{+}
\end{array}\right)
$$

Here, $q_{+}$is such that the boundary conditions minimize $f_{b}$.
After an appropriate rescaling (Bisi et al. 2003), we may actually consider a dimensionless version of the problem, where we are left with only two parameters: a reduced temperature $\theta \in(-\infty, 1)$, and a reduced elastic constant $1 / \lambda^{2}$. The parameter $\lambda>0$ is proportional to $d / \sqrt{L}$ : it accounts for the effects of the elastic constant $L$, and of the distance between the plates $d$. From now on we will work with the reduced free energy

$$
\begin{equation*}
E_{\lambda}(Q)=\int_{-1}^{1}\left(\frac{1}{2 \lambda^{2}}\left|Q^{\prime}\right|^{2}+f(Q)\right) \mathrm{d} x \tag{2}
\end{equation*}
$$

where

$$
\begin{align*}
f(Q) & =\frac{\theta}{6}|Q|^{2}-\frac{2}{3} \operatorname{tr}\left(Q^{3}\right)+\frac{1}{8}|Q|^{4}+c(\theta) \\
& =\frac{\theta}{3}\left(3 q_{1}^{2}+q_{2}^{2}+q_{3}^{2}\right)+4 q_{1}\left(q_{1}^{2}-q_{2}^{2}-q_{3}^{2}\right)+\frac{1}{2}\left(3 q_{1}^{2}+q_{2}^{2}+q_{3}^{2}\right)^{2}+c(\theta) . \tag{3}
\end{align*}
$$

Here, the constant $c(\theta)$ is chosen in such a way that $\min f=0$. Note that this minimum is attained exactly (Majumdar and Zarnescu 2010) at uniaxial $Q$-tensors of the form

$$
Q=6 q_{+}\left(n \otimes n-\frac{1}{3} I\right), \quad n \in \mathbb{S}^{2}, 6 q_{+}=1+\sqrt{1-\theta}
$$

Although we do not emphasize it in the notation, the free energy obviously depends on $\theta$.

The direct method of the calculus of variations applies to the energy functional (2) in the natural space $H^{1}(-1,1)^{3}$. Hence minimizers always exist. They are critical points of the energy, and as such they satisfy the Euler-Lagrange equation

$$
\begin{equation*}
\frac{1}{\lambda^{2}} Q^{\prime \prime}=\frac{\theta}{3} Q-2\left(Q^{2}-\frac{|Q|^{2}}{3} I\right)+\frac{1}{2}|Q|^{2} Q . \tag{4}
\end{equation*}
$$

Solutions of (4) are analytic, and they satisfy the maximum principle (Majumdar and Zarnescu 2010)

$$
\begin{equation*}
|Q| \leq 2 \sqrt{6} q_{+} . \tag{5}
\end{equation*}
$$

In terms of $q_{1}, q_{2}$, and $q_{3}$ defined by (1), the Euler-Lagrange equation (4) becomes the system

$$
\left\{\begin{array}{l}
\frac{1}{\lambda^{2}} q_{1}^{\prime \prime}=\frac{\theta}{3} q_{1}-\frac{2}{3}\left(q_{2}^{2}+q_{3}^{2}-3 q_{1}^{2}\right)+\left(3 q_{1}^{2}+q_{2}^{2}+q_{3}^{2}\right) q_{1}  \tag{6}\\
\frac{1}{\lambda^{2}} q_{2}^{\prime \prime}=\frac{\theta}{3} q_{2}-4 q_{1} q_{2}+\left(3 q_{1}^{2}+q_{2}^{2}+q_{3}^{2}\right) q_{2} \\
\frac{1}{\lambda^{2}} q_{3}^{\prime \prime}=\frac{\theta}{3} q_{3}-4 q_{1} q_{3}+\left(3 q_{1}^{2}+q_{2}^{2}+q_{3}^{2}\right) q_{3}
\end{array}\right.
$$

and the boundary conditions read

$$
\begin{array}{ll}
q_{1}(-1)=q_{+}, & q_{1}(1)=q_{+}, \\
q_{2}(-1)=3 q_{+}, & q_{2}(1)=-3 q_{+},  \tag{7}\\
q_{3}(-1)=0, & q_{3}(1)=0 .
\end{array}
$$

In the sequel we will denote by $\mathscr{H}$ the space of all admissible configurations, i.e., the space of $H^{1}$ configurations satisfying the boundary conditions. Thus, $\mathscr{H}$ is an affine subspace of $H^{1}(-1,1)^{3}$, consisting of all $Q$-tensors of the form (1), which satisfy the boundary conditions (7).

## 3 Eigenvalue Exchange Configurations

Consider the group $G$ defined as the subgroup of $O(3)$ generated by the matrices $S_{y}$ and $S_{z}$ of the orthogonal reflections with respect to the axes $\mathbb{R} \overrightarrow{e_{y}}$ and $\mathbb{R} \overrightarrow{e_{z}}$. Explicitly,

$$
S_{y}=\left(\begin{array}{ccc}
-1 & 0 & 0 \\
0 & 1 & 0 \\
0 & 0 & -1
\end{array}\right) \quad \text { and } \quad S_{z}=\left(\begin{array}{ccc}
-1 & 0 & 0 \\
0 & -1 & 0 \\
0 & 0 & 1
\end{array}\right)
$$

As a subgroup of $O(3), G$ acts naturally on symmetric traceless matrices, and thus on $H^{1}(-1,1)^{3}$, via the following formula:

$$
(R \cdot Q)(x)=R Q(x)^{t} R, \quad R \in G .
$$

One easily sees that the affine subspace $\mathscr{H} \subset H^{1}(-1,1)^{3}$ of admissible configurations is stable under this action: if $Q$ satisfies the boundary conditions (7) then $R \cdot Q$ satisfies them also, for $R \in G$. Thus, $G$ acts on $\mathscr{H}$.

Moreover, the free energy functional $E_{\lambda}$ is invariant under this action:

$$
E_{\lambda}(R \cdot Q)=E_{\lambda}(Q) \quad \forall R \in G, Q \in \mathscr{H} .
$$

Therefore the principle of symmetric criticality (Palais 1979) ensures that critical points among $G$-invariant configurations are critical points of $E_{\lambda}$, that is solutions of the Euler-Lagrange system (6).

More precisely, we denote by $\mathscr{H}^{\text {ee }}$ the affine subspace of $\mathscr{H}$ consisting of all invariant configurations, and by $E_{\lambda}^{\text {ee }}=\left.E_{\lambda}\right|_{\mathscr{C}}$ ee the free energy functional restricted to invariant configurations. It is straightforward to check that

$$
\mathscr{H}^{\text {ee }}=\left\{Q \in \mathscr{H} ; q_{3} \equiv 0\right\}
$$

and the principle of symmetric criticality simply asserts that critical points of $E_{\lambda}^{\mathrm{ee}}$ correspond to solutions of (6) with $q_{3} \equiv 0$. Of course this fact could also be checked by a direct computation.

The elements of $\mathscr{H}^{\text {ee }}$ are the eigenvalue exchange configurations, since $\chi \in \mathscr{H}^{\text {ee }}$ corresponds to $\left(q_{1}, q_{2}\right)$ via

$$
\chi(x)=\left(\begin{array}{ccc}
-2 q_{1}(x) & 0 & 0 \\
0 & q_{1}(x)-q_{2}(x) & 0 \\
0 & 0 & q_{1}(x)+q_{2}(x)
\end{array}\right), \quad x \in[-1,1] .
$$

The free energy of such a $\chi \in \mathscr{H}^{\text {ee }}$ is given by

$$
\begin{align*}
E_{\lambda}^{\mathrm{ee}}(\chi)= & \int_{-1}^{1}\left(\frac{1}{2 \lambda^{2}}\left|\chi^{\prime}\right|^{2}+f(\chi)\right) \mathrm{d} x \\
= & \int_{-1}^{1}\left(\frac{3\left(q_{1}^{\prime}\right)^{2}+\left(q_{2}^{\prime}\right)^{2}}{\lambda^{2}}+\frac{\theta}{3}\left(3 q_{1}^{2}+q_{2}^{2}\right)\right.  \tag{8}\\
& \left.+4 q_{1}\left(q_{1}^{2}-q_{2}^{2}\right)+\frac{1}{2}\left(3 q_{1}^{2}+q_{2}^{2}\right)^{2}+c(\theta)\right) \mathrm{d} x
\end{align*}
$$

and critical points of $E_{\lambda}^{\mathrm{ee}}$ solve the boundary value problem

$$
\begin{cases}\frac{1}{\lambda^{2}} q_{1}^{\prime \prime} & =\frac{\theta}{3} q_{1}-\frac{2}{3}\left(q_{2}^{2}-3 q_{1}^{2}\right)+\left(3 q_{1}^{2}+q_{2}^{2}\right) q_{1}  \tag{9}\\ \frac{1}{\lambda^{2}} q_{2}^{\prime \prime} & =\frac{\theta}{3} q_{2}-4 q_{1} q_{2}+\left(3 q_{1}^{2}+q_{2}^{2}\right) q_{2} \\ q_{1}( \pm 1) & =q_{+} \\ q_{2}( \pm 1) & =\mp 3 q_{+}\end{cases}
$$

Since the direct method of the calculus of variations applies to $E_{\lambda}^{\mathrm{ee}}$, there always exists an eigenvalue exchange minimizer, which is an equilibrium configuration in $\mathscr{H}$. This eigenvalue exchange equilibrium is stable in $\mathscr{H}^{\text {ee }}$, but need not be stable as an equilibrium among all admissible configurations: in principle, symmetry-breaking perturbations may induce a negative second variation of the total free energy $E_{\lambda}$. To study this phenomenon we need to understand the structure of that second variation.

Consider a family $\chi_{\lambda}=\left(q_{1, \lambda}, q_{2, \lambda}\right)$ of eigenvalue exchange configurations. That is, $\chi_{\lambda}$ is a critical point of $E_{\lambda}^{\mathrm{ee}}$, and hence also of $E_{\lambda}$. The principle of symmetric criticality (see Appendix 1) ensures that the orthogonal decomposition

$$
H_{0}^{1}(-1,1)^{3}=H_{\mathrm{sp}} \oplus H_{\mathrm{sb}}=\left\{\left(h_{1}, h_{2}, 0\right)\right\} \oplus\left\{\left(0,0, h_{3}\right)\right\}
$$

corresponding to the decomposition into "symmetry-preserving" perturbations and "symmetry-breaking" perturbations is also orthogonal for the bilinear form $D^{2} E_{\lambda}\left(\chi_{\lambda}\right)$. Namely, for $H \in H_{0}^{1}(-1,1)^{3}$,

$$
\begin{aligned}
D^{2} E_{\lambda}\left(\chi_{\lambda}\right)[H] & =D^{2} E_{\lambda}\left(\chi_{\lambda}\right)\left[\left(h_{1}, h_{2}, 0\right)\right]+D^{2} E_{\lambda}\left(\chi_{\lambda}\right)\left[\left(0,0, h_{3}\right)\right] \\
& =\Phi_{\lambda}\left[h_{1}, h_{2}\right]+\Psi_{\lambda}\left[h_{3}\right]
\end{aligned}
$$

Here $\Phi=\Phi_{\lambda}$ and $\Psi=\Psi_{\lambda}$ are quadratic forms defined on $H_{0}^{1}(-1,1)^{2}$, respectively, $H_{0}^{1}(-1,1)$, by the above equality. Note that $\Phi_{\lambda}$ is nothing else than $D^{2} E_{\lambda}^{\mathrm{ee}}\left(\chi_{\lambda}\right)$, the second variation of restricted free energy. From the computations in Appendix 3, we obtain

$$
\begin{align*}
\Phi\left[h_{1}, h_{2}\right]=\int_{-1}^{1} & \left\{\frac{6\left(h_{1}^{\prime}\right)^{2}+2\left(h_{2}^{\prime}\right)^{2}}{\lambda^{2}}+6\left(\frac{\theta}{3}+2 q_{1}+9 q_{1}^{2}+q_{2}^{2}\right) h_{1}^{2}\right. \\
& \left.+2\left(\frac{\theta}{3}-4 q_{1}+3 q_{1}^{2}+3 q_{2}^{2}\right) h_{2}^{2}+8 q_{2}\left(3 q_{1}-2\right) h_{1} h_{2}\right\} \mathrm{d} x \tag{10}
\end{align*}
$$

and

$$
\begin{equation*}
\Psi\left[h_{3}\right]=\int_{-1}^{1}\left\{\frac{2\left(h_{3}{ }^{\prime}\right)^{2}}{\lambda^{2}}+2\left(\frac{\theta}{3}-4 q_{1}+3 q_{1}^{2}+q_{2}^{2}\right) h_{3}^{2}\right\} \mathrm{d} x . \tag{11}
\end{equation*}
$$

To the quadratic forms $\Phi_{\lambda}$ and $\Psi_{\lambda}$, we may associate bounded linear operators $\mathscr{M}_{\lambda}$ : $H_{0}^{1}(-1,1)^{2} \rightarrow H^{-1}(-1,1)^{2}$ and $\mathscr{L}_{\lambda}: H_{0}^{1}(-1,1) \rightarrow H^{-1}(-1,1)$ such that

$$
\begin{equation*}
\left\langle\mathscr{M}_{\lambda}\left(h_{1}, h_{2}\right),\left(h_{1}, h_{2}\right)\right\rangle=\Phi_{\lambda}\left[h_{1}, h_{2}\right] \quad \text { and }\left\langle\mathscr{L}_{\lambda} h_{3}, h_{3}\right\rangle=\Psi_{\lambda}\left[h_{3}\right] . \tag{12}
\end{equation*}
$$

Of particular interest to us will be the first eigenvalues of these operators, since they measure the local stability of the eigenvalue exchange equilibrium. We will denote the first eigenvalue of $\mathscr{M}_{\lambda}$ (respectively $\mathscr{L}_{\lambda}$ ) by $\nu(\lambda)$ (respectively $\mu(\lambda)$ ). They are given by the following formulas:

$$
\begin{equation*}
v(\lambda)=\inf \frac{\Phi_{\lambda}\left[h_{1}, h_{2}\right]}{\int\left(h_{1}^{2}+h_{2}^{2}\right)}, \quad \mu(\lambda)=\inf \frac{\Psi_{\lambda}[h]}{\int h^{2}} . \tag{13}
\end{equation*}
$$

## 4 The Limits of Very Large and Very Small Cell Thickness

So far, we know that there always exists an eigenvalue exchange solution. However, as the cell thickness grows larger, the numerics in Palffy-Muhoray (1994) and Bisi et al. (2003) predict the existence of a BD solution, that is, a solution of (6) with $q_{3} \neq 0$. In addition, this solution should be approximately uniaxial. In Proposition 4.1, we study the limiting behavior of minimizers as $\lambda$ grows to infinity, and obtain in fact a convergence toward a uniaxial tensor. In particular, the minimizer cannot stay in $\mathscr{H}^{\text {ee }}$, thus for large $\lambda$ there do exist solutions other than the EE minimizer.

Before stating the result, we should remark that, due to the symmetry of the energy functional, any solution with $q_{3} \neq 0$ automatically gives rise to another, distinct solution. Recall indeed from Sect. 3 that $E_{\lambda}$ is $G$-invariant, where $G$ is the subgroup of $O(3)$ generated by the orthogonal reflections $S_{y}$ and $S_{z}$ (with respect to the $y$-axis and to the $z$-axis). For a $Q$-tensor associated to $\left(q_{1}, q_{2}, q_{3}\right)$ via (1), it holds

$$
S_{y} \cdot Q=\left(\begin{array}{ccc}
-2 q_{1} & 0 & 0 \\
0 & q_{1}-q_{2} & -q_{3} \\
0 & -q_{3} & q_{1}+q_{2}
\end{array}\right)
$$

Therefore, if $Q$ is a solution of (6), then the $Q$-tensor with opposite $q_{3}$ is also solution of (6). Moreover, those two solutions $Q$ and $S_{y} \cdot Q$ have same energy. That is why, when studying the limit of minimizers of $E_{\lambda}$ in Proposition 4.1, we will restrict ourselves to $Q$-tensors satisfying, say, $q_{3}(0) \geq 0$, to ensure the uniqueness of the limit.

The limit of a small elastic constant—which corresponds to a large $\lambda$-has already been studied in Majumdar and Zarnescu (2010) in the three-dimensional case, and in Bauman et al. (2012), Canevari (2014), and Golovaty and Montero (2013) in the two-dimensional case. The one-dimensional case considered in the present article is particularly simple and we obtain the following result.

Proposition 4.1 Let $Q_{\lambda}$ be a minimizer of $E_{\lambda}$, with $q_{3}(0) \geq 0$. It holds

$$
Q_{\lambda} \rightarrow Q_{*} \quad \text { in } H^{1}(-1,1)^{3}
$$

as $\lambda$ tends to $+\infty$, where

$$
Q_{*}(x)=6 q_{+}\left(n_{*}(x) \otimes n_{*}(x)-\frac{1}{3} I\right), \quad n_{*}(x)=\left(\begin{array}{c}
0 \\
\cos \left(\frac{\pi}{4}-\frac{\pi}{4} x\right) \\
\sin \left(\frac{\pi}{4}-\frac{\pi}{4} x\right)
\end{array}\right)
$$

Proof One proves, exactly as in Majumdar and Zarnescu (2010, Lemma 3), that there exists a subsequence

$$
Q_{\lambda_{k}} \longrightarrow Q_{*}=6 q_{+}\left(n \otimes n-\frac{1}{3} I\right) \quad \text { in } H^{1}
$$

where $Q_{*}$ minimizes $\int\left|Q^{\prime}\right|^{2}$ among maps in $\mathscr{H}$ which are everywhere of the form $Q=6 q_{+}(n \otimes n-I / 3)$ —that is, maps $Q$ in $\mathscr{H}$ which satisfy $f(Q)=0$ everywhere.

Since $Q_{*}$ is continuous on $(-1,1)$ it follows from Ball and Zarnescu (2011, Lemma 3) that there exists a unique continuous map $n_{*}:(-1,1) \rightarrow \mathbb{S}^{2}$ such that

$$
Q_{*}(x)=6 q_{+}\left(n_{*}(x) \otimes n_{*}(x)-\frac{1}{3} I\right), \quad n_{*}(-1)=\overrightarrow{e_{z}}
$$

Moreover, by Ball and Zarnescu (2011, Lemma 1), the map $n_{*}$ lies in $H^{1}\left(-1,1 ; \mathbb{S}^{2}\right)$.
Since $Q_{*}$ minimizes $\int\left|Q^{\prime}\right|^{2}$, we deduce that $n_{*}$ minimizes $\int\left|n^{\prime}\right|^{2}$ among maps $n \in$ $H^{1}\left(-1,1 ; \mathbb{S}^{2}\right)$ satisfying the same boundary conditions as $n_{*}$. It holds $n_{*}(-1)=\overrightarrow{e_{z}}$, and the boundary conditions on $Q_{*}$ imply that $n_{*}(1)=\alpha \overrightarrow{e_{y}}$ for some $\alpha= \pm 1$. Using the fact that the geodesics on the sphere $\mathbb{S}^{2}$ are arcs of large circles, we obtain

$$
n_{*}(x)=\left(0, \alpha \cos \left(\frac{\pi}{4}-\frac{\pi}{4} x\right), \sin \left(\frac{\pi}{4}-\frac{\pi}{4} x\right)\right)
$$

On the other hand, since the maps $Q_{\lambda}$ satisfy $q_{3, \lambda}(0) \geq 0$, the limiting map $Q_{*}$ must satisfy also $q_{3, *}(0) \geq 0$. Since the above formula for $n_{*}$ implies that

$$
q_{3, *}(0)=6 \alpha q_{+} \cos (\pi / 4) \sin (\pi / 4)=3 \alpha q_{+}
$$

we conclude that $\alpha \geq 0$, and thus $\alpha=1$. In particular, we obtain the formula for $n_{*}$ in the statement of the theorem. Moreover, we have shown that the limit

$$
Q_{*}=\lim Q_{\lambda_{k}}
$$

is uniquely determined, independently of the converging subsequence. Therefore, we do actually have

$$
Q_{\lambda} \longrightarrow Q_{*} \text { in } H^{1}
$$

as $\lambda$ tends to $+\infty$.
Now we turn to studying the case of a very narrow cell. That is, we investigate the limit $\lambda \rightarrow 0$. The numerics in Palffy-Muhoray (1994) and Bisi et al. (2003) predict that for small $\lambda$, there is only one solution, which is an EE configuration. This is indeed the content of the next result, which is in the same spirit as the similar uniqueness result (Bethuel et al. 1994, Theorem VIII.7) in a Ginzburg-Landau setting.

Proposition 4.2 There exists $\lambda_{0}>0$, such that for any $\lambda \in\left(0, \lambda_{0}\right)$, $E_{\lambda}$ admits a unique critical point $\chi_{\lambda} \in \mathscr{H}^{\text {ee }}$.

Proof The uniqueness is a consequence of a more general result, stated as Theorem 8.1 in Appendix 2. The fact that the unique solution belongs to $\mathscr{H}^{\text {ee }}$ is immediate from the considerations in Sect. 3, since there always exists a solution $\chi \in \mathscr{H}^{\text {ee }}$.

Proposition 4.2 provides us with a family of solutions

$$
\left(0, \lambda_{0}\right) \ni \lambda \mapsto \chi_{\lambda} \in \mathscr{H}^{\mathrm{ee}}
$$

The next result gives further properties of this branch of solutions. Recall from (13) the definitions of $\nu(\lambda)$ and $\mu(\lambda): \nu$ is the first eigenvalue of $D^{2} E_{\lambda}^{\mathrm{ee}}\left(\chi_{\lambda}\right)$, and $\mu$ is the first eigenvalue of $D^{2} E_{\lambda}\left(\chi_{\lambda}\right)$ restricted to symmetry-breaking perturbations.

Proposition 4.3 The map $\lambda \mapsto \chi_{\lambda}$ is smooth on $\left(0, \lambda_{0}\right)$ and can be extended uniquely to a smooth map of eigenvalue exchange solutions

$$
\left(0, \lambda_{*}\right) \rightarrow \mathscr{H}^{\mathrm{ee}}, \lambda \mapsto \chi_{\lambda},
$$

where $\lambda_{*} \in\left[\lambda_{0},+\infty\right]$ is determined by the following property:

$$
\begin{equation*}
\left(v(\lambda)>0 \quad \forall \lambda \in\left(0, \lambda_{*}\right)\right) \text { and } \quad\left(\lambda_{*}=+\infty \text { or } v\left(\lambda_{*}\right)=0\right) \tag{14}
\end{equation*}
$$

Moreover, the map $\lambda \mapsto \mu(\lambda)$ is smooth on $\left(0, \lambda_{*}\right)$.
Remark 4.4 It is not clear wether both alternatives in (14) can actually occur. In Theorem 5.1, we show that $\lambda_{*}=+\infty$ provided $\theta=-8$. In general we do not know if $\lambda_{*}$ could be finite.

Proof of Proposition 4.3 In the proof of Theorem 8.1, $\lambda_{0}$ is chosen in such a way that $E_{\lambda}$ is strictly convex around $\chi_{\lambda}$, and in particular $D^{2} E_{\lambda}\left(\chi_{\lambda}\right)$ is positive for $\lambda \in\left(0, \lambda_{0}\right)$. In fact it is straightforward to check (using Poincare's inequality) that the choice of $\lambda_{0}$ in the proof of Theorem 8.1 ensures that $D^{2} E_{\lambda}\left(\chi_{\lambda}\right)$ is positive definite for $\lambda \in\left(0, \lambda_{0}\right)$. In particular, $D^{2} E_{\lambda}^{\mathrm{ee}}\left(\chi_{\lambda}\right)$ is positive definite, or equivalently, $\nu(\lambda)>0$ for $\lambda \in\left(0, \lambda_{0}\right)$.

Therefore $\chi_{\lambda}$ is a non-degenerate critical point, and we may apply the implicit function theorem to the smooth map

$$
\mathscr{F}:(0,+\infty) \times \mathscr{H}^{\mathrm{ee}} \rightarrow H^{-1}(-1,1)^{2},(\lambda, \chi) \mapsto D^{2} E_{\lambda}^{\mathrm{ee}}(\chi),
$$

around a solution $\left(\lambda, \chi_{\lambda}\right)$ of $\mathscr{F}=0$, for $\lambda \in\left(0, \lambda_{0}\right)$. Since this solution is unique, we deduce that $\lambda \mapsto \chi_{\lambda}$ is given by the implicit function theorem and as such, is smooth.

As long as $D^{2} E_{\lambda}\left(\chi_{\lambda}\right)$ stays positive definite, i.e., $v(\lambda)>0$, we may apply the implicit function theorem to smoothly extend the map $\lambda \mapsto \chi_{\lambda}$, until we reach a $\lambda_{*}$ satisfying (14). Note that the extension is unique since for each $\lambda, \chi_{\lambda}$ is a non-degenerate-an thus isolated-critical point of $E_{\lambda}^{\mathrm{ee}}$.

It remains to prove that $\lambda \mapsto \mu(\lambda)$ is a smooth map. Recall that $\mu(\lambda)$ is the first eigenvalue of the bounded linear operator

$$
\begin{align*}
\mathscr{L}_{\lambda}: & H_{0}^{1}(-1,1) \rightarrow H^{-1}(-1,1) \\
& h \mapsto-\frac{2}{\lambda^{2}} h^{\prime \prime}+2\left(\frac{\theta}{3}-4 q_{1, \lambda}+3 q_{1, \lambda}^{2}+q_{2, \lambda}^{2}\right) h, \tag{15}
\end{align*}
$$

where $\left(q_{1, \lambda}, q_{2, \lambda}\right)=\chi_{\lambda}$. From the smoothness of $\lambda \mapsto \chi_{\lambda}$ we deduce easily that $\mathscr{L}_{\lambda}$ depends smoothly on $\lambda$.

Let us fix $\lambda_{0} \in\left(0, \lambda_{*}\right)$. From the theory of Sturm-Liouville operators, we know that $\mu(\lambda)$ is a simple eigenvalue of $\mathscr{L}_{\lambda}$. In fact, in the terminology of Crandall and Rabinowitz (1973, Definition 1.2), $\mu_{0}$ is an $i$-simple eigenvalue of $\mathscr{L}_{\lambda_{0}}$, where $i$ : $H_{0}^{1} \rightarrow H^{-1}$ is the injection operator. Indeed, since $\mathscr{L}_{\lambda_{0}}$ is Fredholm of index 0 and symmetric, if we fix an eigenfunction $h_{0} \in H_{0}^{1}, \int h_{0}^{2}=1$ associated to $\mu_{0}$, then it holds

$$
\operatorname{Ran}\left(\mathscr{L}_{\lambda_{0}}-\mu_{0} i\right)=\left\{f \in H^{-1} ;<f, h_{0}>=0\right\},
$$

so that $i h_{0} \notin \operatorname{Ran}\left(\mathscr{L}_{\lambda_{0}}-\mu_{0} i\right)$ and the $i$-simplicity of $\mu_{0}$ follows easily.
Therefore we may invoke Crandall and Rabinowitz (1973, Lemma 1.3) to obtain the existence of smooth maps $\lambda \mapsto \tilde{\mu}(\lambda), \lambda \mapsto h_{\lambda}$ defined for $\lambda \approx \lambda_{0}$, such that $\tilde{\mu}(\lambda)$ is the unique eigenvalue of $\mathscr{L}_{\lambda_{0}}$ close enough to $\mu_{0}$, and $h_{\lambda}$ a corresponding eigenfunction.

On the other hand, it can be easily checked that $\lambda \mapsto \mu(\lambda)$ is continuous: upper semi-continuity is obvious since $\mu$ is an infimum of continuous functions, and lower semi-continuity follows from the inequalities

$$
\mu\left(\lambda_{0}\right) \leq \mu(\lambda)+\left\|\mathscr{L}_{\lambda_{0}}-\mathscr{L}_{\lambda}\right\|\left\|h_{\lambda}\right\|_{H^{1}} \leq \mu(\lambda)+C\left\|\mathscr{L}_{\lambda_{0}}-\mathscr{L}_{\lambda}\right\|,
$$

where $h_{\lambda} \in H_{0}^{1}$ is a $L^{2}$-normalized eigenfunction associated to $\mu(\lambda)$, and $\lambda$ is close to $\lambda_{0}$. (Note that $\left\|h_{\lambda}\right\|_{H^{1}}$ is bounded since $\left\langle\mathscr{L}_{\lambda} h_{\lambda}, h_{\lambda}\right\rangle$ is bounded.)

Therefore, for $\lambda$ close enough to $\lambda_{0}, \mu(\lambda)$ is close enough to $\mu_{0}$. Hence by the uniqueness in Crandall and Rabinowitz (1973, Lemma 1.3), $\mu(\lambda)$ must coincide with $\tilde{\mu}(\lambda)$. In particular, $\mu$ is smooth.

Although we did not emphasize this dependence in the notations, everything we have done so far depends on the fixed parameter $\theta \in(-\infty, 1)$. In the next section, we choose a special value for this parameter, $\theta=-8$, at which computations are simplified.

## 5 The Special Temperature $\theta=-8$

Throughout the present section, we assume that $\theta=-8$. As already noticed in Bisi et al. (2003), in that case the equation (9) admits a particularly simple solution, and we are able to say a lot more about the branch of solutions $\lambda \mapsto \chi_{\lambda}$ obtained in Proposition 4.3.

First of all, we obtain more information about the maximal value $\lambda_{*}$ of definition of $\chi_{\lambda}$, and about the eigenvalue $\mu(\lambda)$ measuring the stability with respect to symmetrybreaking perturbations. In fact we are going to prove the following theorem, which is the first of two main results in the present section.

Theorem 5.1 Assume that $\theta=-8$. Then $\lambda_{*}=+\infty$. That is, the unique eigenvalue exchange solution $\chi_{\lambda}$ for small $\lambda$ can be extended to a smooth branch of eigenvalue exchange solutions

$$
(0,+\infty) \rightarrow \mathscr{H}^{\mathrm{ee}}, \lambda \mapsto \chi_{\lambda} .
$$

with $\nu(\lambda)>0$ for all $\lambda>0$. Moreover, there exists $\lambda_{c}>0$ such that

$$
\begin{equation*}
\mu(\lambda)>0 \quad \forall \lambda \in\left(0, \lambda_{c}\right), \quad \mu\left(\lambda_{c}\right)=0, \quad \text { and } \mu^{\prime}\left(\lambda_{c}\right)<0 . \tag{16}
\end{equation*}
$$

In fact it holds $\mu^{\prime}(\lambda)<0$ for all $\lambda$, and $\lim _{+\infty} \mu<0$.
In particular, Theorem 5.1 provides a rigorous justification for part of the bifurcation diagram pictured in Fig. 2. Namely, there is a smooth branch of eigenvalue exchange solutions defined for all $\lambda$ and loosing stability at some critical value of $\lambda$.

The next natural step is to investigate what happens at the critical value $\lambda_{c}$, where the branch of eigenvalue exchange solutions looses stability. This is the content of the second main result of the present section. Let $h_{c} \in \operatorname{ker} \mathscr{L}_{\lambda_{c}}$ (a perturbation responsible for the loss of stability at $\lambda_{c}$ ), and denote by $h_{c}^{\perp} \subset H_{0}^{1}(-1,1)^{3}$ the space of perturbations orthogonal to $\left(0,0, h_{c}\right) \in H_{0}^{1}(-1,1)^{3}$.

Theorem 5.2 Assume $\theta=-8$. There exist $\delta, \varepsilon>0$ and a neighborhood $\mathfrak{A}$ of $\chi_{\lambda_{c}}$ in $\mathscr{H}$, such that the solutions of

$$
\begin{equation*}
D E_{\lambda}(Q)=0, \quad(\lambda, Q) \in\left(\lambda_{c}-\delta, \lambda_{c}+\delta\right) \times \mathfrak{A}, \tag{17}
\end{equation*}
$$

are exactly

$$
Q=\chi_{\lambda} \text { or }\left\{\begin{array}{l}
\lambda=\lambda(t)  \tag{18}\\
Q=\chi_{\lambda_{c}}+t\left(0,0, h_{c}\right)+t^{2} H_{t},
\end{array} \quad \text { for some } t \in(-\varepsilon, \varepsilon)\right.
$$

where $\lambda(t) \in\left(\lambda_{c}-\delta, \lambda_{c}+\delta\right)$ and $H_{t} \in h_{c}^{\perp}$ are smooth functions of $t \in(-\varepsilon, \varepsilon)$. Moreover, the following symmetry properties are satisfied:

$$
\begin{equation*}
\lambda(-t)=\lambda(t), \quad \text { and } \quad h_{1,-t}=h_{1, t}, h_{2,-t}=h_{2, t}, h_{3,-t}=-h_{3, t}, \tag{19}
\end{equation*}
$$

where $H_{t}$ is identified with $\left(h_{1, t}, h_{2, t}, h_{3, t}\right)$ via (1).
The rest of the section will be devoted to the proofs of Theorems 5.1 and 5.2 , which we decompose into several intermediate results.
5.1 The Proof of Theorem 5.1

We start by proving that the eigenvalue exchange solution branch $\chi_{\lambda}$ obtained in Proposition 4.3 has constant $q_{1}$, and can be extended to all $\lambda>0$. In particular we obtain the first part of Theorem 5.1.

Proposition 5.3 Assume $\theta=-8$. Then $\lambda_{*}=+\infty$, and for every $\lambda \in(0,+\infty)$, $\chi_{\lambda}=\left(2 / 3, q_{2, \lambda}\right)$, where $q_{2}=q_{2, \lambda}$ solves

$$
\left\{\begin{array}{l}
\frac{1}{\lambda^{2}} q_{2}^{\prime \prime}=\left(q_{2}^{2}-4\right) q_{2}  \tag{20}\\
q_{2}(-1)=2, q_{2}(1)=-2
\end{array}\right.
$$

Proof When the value of the reduced temperature $\theta$ is set to $\theta=-8$, then $q_{+}=2 / 3$. Let us define $\tilde{q}_{1}=q_{1}-2 / 3$. For $\tilde{q}_{1}$, the boundary conditions become $\tilde{q}_{1}( \pm 1)=0$. The boundary conditions for $q_{2}$ are $q_{2}( \pm 1)=\mp 2$.

In terms of $\tilde{q}_{1}$, the bulk energy density-for the eigenvalue exchange solution (that is, with $q_{3}=0$ )-reads

$$
\begin{equation*}
f\left(q_{1}, q_{2}\right)=16 \tilde{q}_{1}^{2}\left(\frac{3}{2}+\tilde{q}_{1}\right)+\frac{1}{2}\left(q_{2}^{2}-4+3 \tilde{q}_{1}^{2}\right)^{2} \tag{21}
\end{equation*}
$$

and the Euler-Lagrange equations become

$$
\begin{align*}
\frac{1}{\lambda^{2}} \tilde{q}_{1}^{\prime \prime} & =\left(4+8 \tilde{q}_{1}+3 \tilde{q}_{1}^{2}+q_{2}^{2}\right) \tilde{q}_{1} \\
\frac{1}{\lambda^{2}} q_{2}^{\prime \prime} & =\left(q_{2}^{2}-4+3 \tilde{q}_{1}^{2}\right) q_{2} \tag{22}
\end{align*}
$$

Therefore, there exists a solution with $\tilde{q}_{1} \equiv 0$, i.e., $q_{1} \equiv 2 / 3$. Indeed, a constant $\tilde{q}_{1}$ solves the first equation (for any $q_{2}$ ), and the corresponding $q_{2}$ is obtained by minimizing the energy $E_{\lambda}^{\text {ee }}$ in which $q_{1}$ is taken to be constant. That is, $q_{2}$ minimizes

$$
\begin{equation*}
I_{\lambda}\left(q_{2}\right)=\int_{-1}^{1}\left(\frac{1}{\lambda^{2}}\left(q_{2}^{\prime}\right)^{2}+\frac{1}{2}\left(q_{2}^{2}-4\right)^{2}\right) \mathrm{d} x \tag{23}
\end{equation*}
$$

Hence $q_{2}$ solves (20). From Lemma 5.4, we know that (20) actually admits a unique solution. Hence we may define for all $\lambda>0$, without ambiguity, the EE solution

$$
\tilde{\chi}_{\lambda}:=\left(2 / 3, q_{2, \lambda}\right),
$$

where $q_{2, \lambda}$ solves (20).
The uniqueness proven in Proposition 4.2 ensures that $\chi_{\lambda}=\tilde{\chi}_{\lambda}$ for $\lambda \in\left(0, \lambda_{0}\right)$. On the other hand, Lemma 5.5 ensures that $\tilde{\chi}_{\lambda}$ is a smooth extension of $\chi_{\lambda}$ satisfying $\nu(\lambda)>0$ for all $\lambda>0$. Therefore we conclude, by the uniqueness in Proposition 4.3, that $\lambda_{*}=+\infty$ and $\chi_{\lambda}=\left(2 / 3, q_{2, \lambda}\right)$.

In the proof of Proposition 5.3, we made use of two lemmas, Lemma 5.4 and Lemma 5.5, that we are going to prove next. The first one gives properties of the boundary value problem (20) satisfied by $q_{2, \lambda}$.

Lemma 5.4 The boundary value problem (20) has a unique solution, which is odd and decreasing.

Proof Very similar results are classical in the study of reaction-diffusion equations [see for instance Fife (1979, Sect. 4.3)]. Since the present case is particularly simple, we nevertheless give a complete proof here. Recall that the existence of a solution follows directly from minimizing the energy $I_{\lambda}$ defined in (23).

We start by proving the bounds

$$
\begin{equation*}
-2 \leq q_{2} \leq 2 \tag{24}
\end{equation*}
$$

Assume that $q_{2}^{2}$ attains its maximum in $(-1,1)$. Then, at a point where the maximum is attained, it holds

$$
0 \geq \frac{1}{2 \lambda^{2}}\left(q_{2}^{2}\right)^{\prime \prime} \geq \frac{1}{\lambda^{2}} q_{2}^{\prime \prime} q_{2}=\left(q_{2}^{2}-4\right) q_{2}^{2}
$$

so that $q_{2}^{2} \leq 4$. Since this bound is satisfied (with equality) on the boundary, (24) is proved.

Multiplying (20) by $q_{2}{ }^{\prime}$, we obtain the first integral

$$
\begin{equation*}
\left[\frac{1}{2 \lambda^{2}}\left(q_{2}^{\prime}\right)^{2}\right]^{\prime}=\left[\frac{1}{4}\left(q_{2}^{2}-4\right)^{2}\right]^{\prime} \tag{25}
\end{equation*}
$$

Integrating (25), we obtain

$$
\begin{equation*}
\frac{1}{2 \lambda^{2}}\left(q_{2}^{\prime}\right)^{2}=\frac{1}{4}\left(q_{2}^{2}-4\right)^{2}+\frac{1}{2 \lambda^{2}} q_{2}^{\prime}(-1)^{2} . \tag{26}
\end{equation*}
$$

Since $q_{2}{ }^{\prime}(-1) \neq 0$ (otherwise $q_{2}$ would satisfy the same Cauchy problem at -1 as the constant solution), it follows in particular that $q_{2}{ }^{\prime}$ does not vanish. On the other hand, the bounds (24) ensure that $q_{2}{ }^{\prime}(-1)$ is negative. Therefore $q_{2}{ }^{\prime}$ must stay negative:

$$
\begin{equation*}
q_{2}^{\prime}<0, \tag{27}
\end{equation*}
$$

hence every solution of (20) is decreasing.
Now we prove that (20) has a unique solution. Assume $\underline{q}_{2}$ and $\bar{q}_{2}$ are distinct solution. Then, they must have distinct derivatives at -1 (otherwise they would satisfy the same Cauchy problem). Say

$$
\begin{equation*}
\underline{q}_{2}^{\prime}(-1)<\bar{q}_{2}^{\prime}(-1)<0 . \tag{28}
\end{equation*}
$$

Since $\underline{q}_{2}$ and $\bar{q}_{2}$ take the same value at 1 , we may consider

$$
x_{0}=\min \left\{x>-1 ; \underline{q}_{2}(x)=\bar{q}_{2}(x)\right\} \in(-1,1] .
$$

At this point $x_{0}, \underline{q}_{2}$ and $\bar{q}_{2}$ must have distinct derivatives, and since $\underline{q}_{2}<\bar{q}_{2}$ in $\left(-1, x_{0}\right)$, it holds

$$
\begin{equation*}
\bar{q}_{2}^{\prime}\left(x_{0}\right)<\underline{q}_{2}^{\prime}\left(x_{0}\right)<0 \tag{29}
\end{equation*}
$$

From (28) and (29) we deduce that

$$
\bar{q}_{2}^{\prime}(-1)^{2}-\bar{q}_{2}^{\prime}\left(x_{0}\right)^{2}<\underline{q}_{2}^{\prime}(-1)^{2}-\underline{q}_{2}^{\prime}\left(x_{0}\right)^{2}
$$

which is obviously incompatible with the facts that $\underline{q}_{2}$ and $\bar{q}_{2}$ satisfy (25) and coincide at -1 and $x_{0}$. Therefore (20) has a unique solution.

Next we prove that $q_{2}$ satisfying (20) must be odd. This is a direct consequence of uniqueness: the functions $q_{2}(x)$ and $-q_{2}(-x)$ are both solutions of the boundary problem (20), therefore they must coincide.

Now we turn to the proof of the second lemma used in the proof of Proposition 5.3, in which we show that the EE solution with constant $q_{1}$ is non-degenerately stable in $\mathscr{H}^{\mathrm{ee}}$.

Lemma 5.5 Assume $\theta=-8$. Let $q_{2, \lambda}$ be the unique solution of (20), and $\chi_{\lambda}:=$ $\left(2 / 3, q_{2, \lambda}\right) \in \mathscr{H}^{\text {ee }}$. Then $\nu(\lambda)$, defined as in (13), satisfies

$$
\begin{equation*}
\nu(\lambda)>0 \quad \forall \lambda>0 . \tag{30}
\end{equation*}
$$

As a consequence, $\lambda \mapsto q_{2, \lambda}$ is smooth.
Proof First note that the smoothness of $\lambda \mapsto q_{2, \lambda}$ follows from (30). Indeed, (30) implies that $D^{2} I_{\lambda}\left(q_{2, \lambda}\right)$ is invertible, so that near $q_{2, \lambda}$, a solution $q_{2}$ of $D I_{\lambda}\left(q_{2}\right)=0$ depending smoothly on $\lambda$ may be obtained by the implicit function theorem. On the other hand, the uniqueness proven in Lemma 5.4 implies that $q_{2, \lambda}$ coincide with this smooth solution.

Now we turn to the proof of (30). Recall that $v(\lambda)$ is the first eigenvalue of the quadratic form $\Phi_{\lambda}=D^{2} E_{\lambda}^{\mathrm{ee}}\left(\chi_{\lambda}\right)$. Since $\theta=-8$ and $q_{1} \equiv 2 / 3$, it holds

$$
\begin{aligned}
\Phi_{\lambda}\left[h_{1}, h_{2}\right]= & \int_{-1}^{1}\left\{\frac{6}{\lambda^{2}}\left(h_{1}^{\prime}\right)^{2}+6\left(\frac{8}{3}+q_{2, \lambda}^{2}\right) h_{1}^{2}\right\} \mathrm{d} x \\
& +\int_{-1}^{1}\left\{\frac{2}{\lambda^{2}}\left(h_{2}^{\prime}\right)^{2}+2\left(3 q_{2, \lambda}^{2}-4\right) h_{2}^{2}\right\} \mathrm{d} x .
\end{aligned}
$$

That is, $\Phi_{\lambda}$ decomposes into a quadratic form in $h_{1}$, which is obviously positive definite, and a quadratic form in $h_{2}$, which is nothing else than $D^{2} I_{\lambda}\left(q_{2, \lambda}\right)$. Therefore, to prove (30) we only need to show that $D^{2} I_{\lambda}\left(q_{2, \lambda}\right)$ is positive definite.

Let us define

$$
\begin{align*}
\eta(\lambda) & :=\inf _{h \in H_{0}^{1}(-1,1), \int h^{2}=1} D^{2} I_{\lambda}\left(q_{2, \lambda}\right)[h] \\
& =\inf _{h \in H_{0}^{1}(-1,1) \int h^{2}=1} \int\left\{\frac{2}{\lambda^{2}}\left(h^{\prime}\right)^{2}+2\left(3 q_{2, \lambda}^{2}-4\right) h^{2}\right\} \mathrm{d} x . \tag{31}
\end{align*}
$$

We need to prove that $\eta(\lambda)>0$ for every $\lambda>0$. Since $q_{2, \lambda}$ minimizes $I_{\lambda}$, it clearly holds $\eta(\lambda) \geq 0$. To prove that $\eta(\lambda)$ cannot vanish, we are going to establish that $\lambda \mapsto \eta(\lambda)$ is decreasing.

To this end, we remark that after a rescaling it holds

$$
\eta(\lambda)=\inf _{h \in H_{0}^{1}(-\lambda, \lambda), \int h^{2}=1} \int_{-\lambda}^{\lambda}\left\{2\left(h^{\prime}\right)^{2}+2\left(3 \bar{q}_{2, \lambda}^{2}-4\right) h^{2}\right\} \mathrm{d} y
$$

where $\bar{q}_{2, \lambda}$ is the rescaled map defined by

$$
\begin{equation*}
\bar{q}_{2, \lambda}(y)=q_{2, \lambda}(y / \lambda), \tag{32}
\end{equation*}
$$

and extended to the whole real line by putting $\bar{q}_{2, \lambda}=2$ in $(-\infty,-\lambda)$ and $\bar{q}_{2, \lambda}=-2$ in $(\lambda,+\infty)$. In Lemma 5.6, we show that $\bar{q}_{2, \lambda}(y)$ is a monotone function of $\lambda$.

Using Lemma 5.6 , we prove that $\eta(\lambda)$ is decreasing: let $\lambda^{\prime}>\lambda$ and consider a map $h_{\lambda} \in H_{0}^{1}(-\lambda, \lambda)$ at which the infimum defining $\eta(\lambda)$ is attained. Then, $h_{\lambda}$ is admissible in the infimum defining $\eta\left(\lambda^{\prime}\right)$, and we obtain $\eta\left(\lambda^{\prime}\right)<\eta(\lambda)$, since $\bar{q}_{2, \lambda^{\prime}}^{2} \leq \bar{q}_{2, \lambda}^{2}$, with strict inequality on $(-\lambda, 0) \cup(0, \lambda)$. The latter fact follows from Lemma 5.6 and the fact that $q_{2, \lambda}$ is odd.

We may now complete the proof of Lemma 5.5: since $\eta(\lambda)$ decreases, and in addition $\eta(\lambda) \geq 0$ for all $\lambda$, we must have $\eta(\lambda)>0$ for any $\lambda$.

In the following lemma, we prove the monotonicity of $\lambda \mapsto \bar{q}_{2, \lambda}$.
Lemma 5.6 For any $y>0,(0, y) \ni \lambda \mapsto q_{2, \lambda}(y / \lambda)=\bar{q}_{2, \lambda}(y)$ is increasing.
Proof The rescaled map $\bar{q}_{2, \lambda}$ minimizes the energy functional

$$
\tilde{I}_{\lambda}\left(\bar{q}_{2}\right)=\int_{0}^{\lambda}\left[\left(\bar{q}_{2}^{\prime}\right)^{2}+\frac{1}{2}\left(\bar{q}_{2}^{2}-4\right)^{2}\right] \mathrm{d} y
$$

subject to the boundary conditions $\bar{q}_{2}(0)=0, \bar{q}_{2}(\lambda)=-2$. Note that we were able to restrict the integral to the positive half-line since $q_{2}$ is odd.

Let $\lambda^{\prime}>\lambda>0$. Consider the respective minimizers $\bar{q}_{2, \lambda^{\prime}}$ and $\bar{q}_{2, \lambda}$, and assume that it does not hold

$$
\bar{q}_{2, \lambda^{\prime}}(y)>\bar{q}_{2, \lambda}(y) \quad \forall y \in\left(0, \lambda^{\prime}\right) .
$$

Then $\bar{q}_{2, \lambda^{\prime}}\left(y_{0}\right)=\bar{q}_{2, \lambda}\left(y_{0}\right)$ for some $y_{0} \in(0, \lambda)$, since in $\left(\lambda, \lambda^{\prime}\right)$ it does hold $\bar{q}_{2, \lambda^{\prime}}>$ $\bar{q}_{2, \lambda}$.

Thus, the maps

$$
\tilde{q}_{2, \lambda}=\bar{q}_{2, \lambda^{\prime}} \mathbf{1}_{y \leq y_{0}}+\bar{q}_{2, \lambda} \mathbf{1}_{y \geq y_{0}}, \quad \tilde{q}_{2, \lambda^{\prime}}=\bar{q}_{2, \lambda} \mathbf{1}_{y \leq y_{0}}+\bar{q}_{2, \lambda^{\prime}} \mathbf{1}_{y \geq y_{0}}
$$

belong to $H^{1}(0, \lambda)$, respectively, $H^{1}\left(0, \lambda^{\prime}\right)$. We claim that $\tilde{q}_{2, \lambda}$ minimizes $\tilde{I}_{\lambda}$. Assume indeed that $\tilde{q}_{2, \lambda}$ has strictly higher energy than $\bar{q}_{2, \lambda}$ : then $\tilde{q}_{2, \lambda^{\prime}}$ would have strictly lower energy than $\bar{q}_{2, \lambda^{\prime}}$, which is absurd. In particular, $\tilde{q}_{2, \lambda}$ is analytical (as a minimizer of $\tilde{I}_{\lambda}$ ), which is possible only if $q_{2, \lambda^{\prime}}$ and $q_{2, \lambda}$ coincide. But then the analytical function $q_{2, \lambda^{\prime}}$ would be constant on $\left(\lambda, \lambda^{\prime}\right)$, and we obtain a contradiction.

So far we have proven the first part of Theorem 5.1, about the extension of $\chi_{\lambda}$ until $\lambda=+\infty$. Now we turn to proving the second part, about the behavior of $\mu(\lambda)$. We split this second part into Propositions 5.7 and 5.8. We start by showing that $\mu(\lambda)$ decreases, with non-vanishing derivative.

Proposition 5.7 Assume $\theta=-8$. Then it holds

$$
\mu^{\prime}(\lambda)<0,
$$

for all $\lambda>0$.
Proof The fact that $\mu(\lambda)$ decreases can be obtained quite easily as a consequence of Lemma 5.6. The fact that its derivative does not vanish, however, is not immediate. Our proof is very similar to the proof of Alama et al. (2012, Proposition 5.18).

First we show that

$$
\frac{\partial}{\partial \lambda}\left[\bar{q}_{2, \lambda}(x)\right]>0 \text { for } x \in(0, \lambda] .
$$

Consider the smooth map

$$
\phi:[0,+\infty) \times \mathbb{R} \rightarrow \mathbb{R},(x, \alpha) \mapsto \phi(x, \alpha),
$$

defined as the solution of the Cauchy problem

$$
\left\{\begin{array}{l}
\phi_{x x}=\left(\phi^{2}-4\right) \phi, \\
\phi(0, \alpha)=0, \phi_{x}(0, \alpha)=\alpha .
\end{array}\right.
$$

Clearly, for any $\lambda>0$, and for $x \in(0, \lambda]$,

$$
\bar{q}_{2, \lambda}(x)=\phi\left(x, \alpha_{\lambda}\right) \quad \text { with } \alpha_{\lambda}=\bar{q}_{2, \lambda^{\prime}}(0) .
$$

Notice that $\alpha_{\lambda}$ solves

$$
\phi\left(\lambda, \alpha_{\lambda}\right)=-2 .
$$

We claim that, for any $x \in(0, \lambda], \partial_{\alpha} \phi\left(x, \alpha_{\lambda}\right)>0$. In fact, let $h(x)=\partial_{\alpha} \phi\left(x, \alpha_{\lambda}\right)$. The function $h$ solves

$$
\left\{\begin{array}{l}
h^{\prime \prime}=\left(3 \bar{q}_{2, \lambda}^{2}-4\right) h, \\
h(0)=0, \quad h^{\prime}(0)=1
\end{array}\right.
$$

Assume that $h\left(x_{0}\right)=0$ for some $x_{0} \in(0, \lambda]$. Then, $h \mathbf{1}_{\left(0, x_{0}\right)}$ would be an admissible test function in the variational problem defining $\eta(\lambda)$, and we would obtain $\eta(\lambda)=0$, which is not true. Recall that $\eta(\lambda)$ was defined in (31), as the first eigenvalue of the second variation of $I_{\lambda}$, and that we have shown in Lemma 5.5 that $\eta(\lambda)>0$ for every $\lambda$.

In particular, $\partial_{\alpha} \phi\left(\lambda, \alpha_{\lambda}\right)>0$ and we can apply the implicit function theorem to obtain a smooth map $\lambda \mapsto \alpha(\lambda)$ such that

$$
\phi(\lambda, \alpha(\lambda))=-2 .
$$

Then, $x \mapsto \phi(x, \alpha(\lambda))$ solves the same boundary problem as $\bar{q}_{2, \lambda}$. By the uniqueness proven in Lemma 5.4, those two functions must be equal, and in particular we have $\alpha_{\lambda}=\alpha(\lambda)$. Moreover, differentiating the equation satisfied by $\alpha(\lambda)$, we obtain

$$
\alpha^{\prime}(\lambda)=-\frac{\partial_{x} \phi\left(\lambda, \alpha_{\lambda}\right)}{\partial_{\alpha} \phi\left(\lambda, \alpha_{\lambda}\right)}=-\frac{\bar{q}_{2, \lambda^{\prime}}(\lambda)}{\partial_{\alpha} \phi\left(\lambda, \alpha_{\lambda}\right)}>0 .
$$

In fact, the bounds (24) ensure that $\bar{q}_{2, \lambda^{\prime}}(\lambda) \leq 0$, and equality cannot occur, else $\bar{q}_{2, \lambda}$ would satisfy the same Cauchy problem as the constant map $q \equiv-2$.

Thus, we have

$$
\frac{\partial}{\partial \lambda}\left[\bar{q}_{2, \lambda}(x)\right]=\alpha^{\prime}(\lambda) \partial_{\alpha} \phi\left(x, \alpha_{\lambda}\right)>0 \text { for } x \in(0, \lambda] .
$$

Let $\lambda_{1}>\lambda_{0}>0$. Using the facts that $(x, \lambda) \mapsto \bar{q}_{2, \lambda}(x)$ is smooth, that $\bar{q}_{2, \lambda}<0$ on $(0,+\infty)$, and that $\bar{q}_{2, \lambda}{ }^{\prime}(0)<0$ (else $\bar{q}_{2, \lambda}$ would coincide with the constant solution $q \equiv 0$ ), we obtain

$$
\bar{q}_{2, \lambda}(x) \leq-c x \quad \forall x \in[0, \lambda], \lambda \in\left[\lambda_{0}, \lambda_{1}\right],
$$

for some constant $c>0$. Similarly, we have

$$
\partial_{\lambda}\left[\bar{q}_{2, \lambda}(x)\right] \geq c^{\prime} x \quad \forall x \in[0, \lambda], \lambda \in\left[\lambda_{0}, \lambda_{1}\right] .
$$

Therefore, we deduce from the mean value theorem the existence of a constant $C>0$ such that

$$
\begin{equation*}
\bar{q}_{2, \lambda}^{2}(x)-\bar{q}_{2, \lambda_{0}}^{2}(x) \leq-C\left(\lambda-\lambda_{0}\right) x^{2} \quad \forall x \in\left(0, \lambda_{0}\right), \lambda \in\left[\lambda_{0}, \lambda_{1}\right] . \tag{33}
\end{equation*}
$$

Note that, since $\bar{q}_{2, \lambda}$ is odd, estimate (33) holds also for all $x \in\left(-\lambda_{0}, \lambda_{0}\right)$.
We remark that, since $\theta=-8$ and $q_{1} \equiv 2 / 3$, formula (13) for $\mu(\lambda)$ simplifies to

$$
\begin{aligned}
\mu(\lambda) & =2 \inf _{h \in H_{0}^{1}(-1,1), \int h^{2}=1} \int_{-1}^{1}\left(\frac{1}{\lambda^{2}}\left(h^{\prime}\right)^{2}+\left(q_{2, \lambda}^{2}-4\right) h^{2}\right) \mathrm{d} x \\
& =2 \inf _{h \in H_{0}^{1}(-\lambda, \lambda) \int h^{2}=1} \int_{-\lambda}^{\lambda}\left(\left(h^{\prime}\right)^{2}+\left(\bar{q}_{2, \lambda}^{2}-4\right) h^{2}\right) \mathrm{d} y .
\end{aligned}
$$

Let $h_{0} \in H_{0}^{1}(-\lambda, \lambda), \int h_{0}^{2}=1$, be a function at which the infimum defining $\mu\left(\lambda_{0}\right)$ is attained. Using the estimate (33), we compute, for $\lambda \in\left[\lambda_{0}, \lambda_{1}\right]$ :

$$
\begin{aligned}
\mu(\lambda) & =2 \inf _{h \in H_{0}^{1}(-\lambda, \lambda), \int h^{2}=1} \int_{-\lambda}^{\lambda}\left[\left(h^{\prime}\right)^{2}+\left(\bar{q}_{2, \lambda}^{2}-4\right) h^{2}\right] \mathrm{d} x \\
& \leq 2 \int_{-\lambda_{0}}^{\lambda_{0}}\left[\left(h_{0}^{\prime}\right)^{2}+\left(\bar{q}_{2, \lambda}^{2}-4\right) h_{0}^{2}\right] \mathrm{d} x \\
& \leq \mu\left(\lambda_{0}\right)-2 C\left(\lambda-\lambda_{0}\right) \int h_{0}^{2} x^{2} \mathrm{~d} x,
\end{aligned}
$$

so that $\mu^{\prime}\left(\lambda_{0}\right)<0$.
To complete the proof of Theorem 5.1, it remains to show that, for large $\lambda$, the eigenvalue exchange solution is unstable with respect to symmetry breaking perturbations. This is the content of the next result.

Proposition 5.8 Assume $\theta=-8$. Then it holds

$$
\lim _{\lambda \rightarrow+\infty} \mu(\lambda)<0
$$

Proof We start by studying the limit of the rescaled map $\bar{q}_{2, \lambda}(y)=q_{2, \lambda}(y / \lambda)$ (extended to $(-\infty,+\infty)$ by $\bar{q}_{2, \lambda} \equiv \mp 2$ near $\pm \infty$ ). This rescaled map $\bar{q}_{2, \lambda}$ minimizes the integral

$$
J_{\lambda}\left(\bar{q}_{2}\right)=\int_{-\lambda}^{\lambda}\left(\left(\bar{q}_{2}^{\prime}\right)^{2}+\frac{1}{2}\left(\bar{q}_{2}^{2}-4\right)^{2}\right) \mathrm{d} y
$$

subject to the boundary conditions $\bar{q}_{2}( \pm \lambda)=\mp 2$. For $\lambda^{\prime}>\lambda, \bar{q}_{2, \lambda}$ is admissible in $J_{\lambda^{\prime}}$. Therefore we deduce that

$$
\lambda \mapsto J_{\lambda}\left(\bar{q}_{2, \lambda}\right) \text { is non-increasing. }
$$

In particular, it holds

$$
\int_{-\infty}^{+\infty}\left(\left(\bar{q}_{2, \lambda}^{\prime}\right)^{2}+\frac{1}{2}\left(\bar{q}_{2, \lambda}^{2}-4\right)^{2}\right) \mathrm{d} x \leq C
$$

and $\left(\bar{q}_{2, \lambda}\right)_{\lambda>0}$ is bounded in $H_{l o c}^{1}(\mathbb{R})$, so that we may extract a weakly converging subsequence. On the other hand, we know from Lemma 5.6 that $\bar{q}_{2, \lambda}(y)$ is a monotonic function of $\lambda$, so that the whole sequence converges pointwise. Therefore the weak $H_{l o c}^{1}$ limit is unique and we do not need to take a subsequence: there exists $\bar{q}_{2, *} \in H_{l o c}^{1}(\mathbb{R})$ such that $\bar{q}_{2, \lambda}$ converges to $\bar{q}_{2, *}$ as $\lambda \rightarrow+\infty$, on every compact interval, $H^{1}$-weakly and uniformly. Using the differential equation satisfied by $\bar{q}_{2, \lambda}$, we see that the second derivatives converge uniformly on every compact interval, so that we actually obtain
convergence in $C_{\text {loc }}^{2}(\mathbb{R})$. In particular, the rescaled limiting map $\bar{q}_{2, *} \in C^{2}(\mathbb{R})$ solves the equation

$$
\begin{equation*}
\bar{q}_{2, *}^{\prime \prime}=\left(\bar{q}_{2, *}^{2}-4\right) \bar{q}_{2, *} . \tag{34}
\end{equation*}
$$

Moreover, using Fatou's lemma, we find that the map $\bar{q}_{2, *}$ has finite energy:

$$
\begin{equation*}
\int_{-\infty}^{+\infty}\left(\left(\bar{q}_{2, *^{\prime}}\right)^{2}+\frac{1}{2}\left(\bar{q}_{2, *}^{2}-4\right)^{2}\right) \mathrm{d} y<+\infty \tag{35}
\end{equation*}
$$

Since $\bar{q}_{2, *}$ is obviously odd and non-increasing, the finite energy property implies that it satisfies the boundary conditions

$$
\bar{q}_{2, *}(-\infty)=2, \quad \bar{q}_{2, *}(+\infty)=-2 .
$$

Recall that, since $\theta=-8$ and $q_{1} \equiv 2 / 3$,

$$
\mu(\lambda)=2 \inf _{h \in H_{0}^{1}(-\lambda, \lambda), \int h^{2}=1} \int_{-\lambda}^{\lambda}\left(\left(h^{\prime}\right)^{2}+\left(\bar{q}_{2, \lambda}^{2}-4\right) h^{2}\right) \mathrm{d} x .
$$

We claim that the convergence of $\bar{q}_{2, \lambda}$ toward $\bar{q}_{2, *}$ implies that

$$
\begin{equation*}
\lim _{\lambda \rightarrow+\infty} \mu(\lambda)=2 \inf _{h \in H^{1}(\mathbb{R}), \int h^{2}=1} \int_{-\infty}^{+\infty}\left(\left(h^{\prime}\right)^{2}+\left(\bar{q}_{2, *}^{2}-4\right) h^{2}\right) \mathrm{d} x . \tag{36}
\end{equation*}
$$

Indeed, for any $\varepsilon>0$, we may find $h_{0} \in C_{c}^{\infty}(\mathbb{R})$ such that

$$
\int_{-\infty}^{+\infty}\left(\left(h_{0}^{\prime}\right)^{2}+\left(\bar{q}_{2, *}^{2}-4\right) h_{0}^{2}\right) \mathrm{d} x \leq m+\varepsilon
$$

where $m$ denotes the infimum in the right-hand side of (36). Choose $\Lambda>0$ such that $\operatorname{supp} h_{0} \subset[-\Lambda, \Lambda]$. Then, for any $\lambda \geq \Lambda$, it holds

$$
\mu(\lambda) \leq 2 \int_{-\infty}^{+\infty}\left(\left(h_{0}{ }^{\prime}\right)^{2}+\left(\bar{q}_{2, \lambda}^{2}-4\right) h_{0}^{2}\right) \mathrm{d} x .
$$

Since $\bar{q}_{2, \lambda}$ converges uniformly to $\bar{q}_{2, *}$ on supp $h_{0}$, we may pass to the limit in the last inequality, and deduce that

$$
\lim _{\lambda \rightarrow+\infty} \mu(\lambda) \leq 2 m+2 \varepsilon,
$$

which proves (36) since $\varepsilon$ is arbitrary.
In view of (36), to conclude the proof we need to find a function $h \in H^{1}(\mathbb{R}), h \neq 0$, such that

$$
\int_{-\infty}^{+\infty}\left(\left(h^{\prime}\right)^{2}+\left(\bar{q}_{2, *}^{2}-4\right) h^{2}\right) \mathrm{d} x<0
$$

We claim that $h=\bar{q}_{2, *}{ }^{\prime}$ is a suitable choice. The fact that $h \neq 0$ is clear in view of the boundary conditions satisfied by $\bar{q}_{2, *}$. The fact that $h \in H_{0}^{1}(\mathbb{R})$ follows from the finite energy property (35). Indeed, (35) clearly implies that $h \in L^{2}(\mathbb{R})$, and also that $\left(\bar{q}_{2, *}^{2}-4\right) \in L^{2}$, so that

$$
h^{\prime}=\left(\bar{q}_{2, *}^{2}-4\right) \bar{q}_{2, *} \in L^{2}(\mathbb{R})
$$

since $\bar{q}_{2, *} \in L^{\infty}$.
Moreover, differentiating the equation satisfied by $\bar{q}_{2, *}$, we obtain

$$
h^{\prime \prime}=\left(3 \bar{q}_{2, *}^{2}-4\right) h,
$$

so that

$$
\int_{-\infty}^{+\infty}\left(\left(h^{\prime}\right)^{2}+\left(\bar{q}_{2, *}^{2}-4\right) h^{2}\right) \mathrm{d} x=-2 \int_{-\infty}^{+\infty} \bar{q}_{2, *}^{2} h^{2} \mathrm{~d} x<0,
$$

and the proof is complete.
Now Theorem 5.1 is obtained directly by putting together the Propositions 5.3, 5.7, and 5.8 above.

### 5.2 The Proof of Theorem 5.2

We define the map

$$
\mathscr{G}:(0,+\infty) \times H_{0}^{1}(-1,1)^{3} \longrightarrow H^{-1}(-1,1)^{3},
$$

defined by

$$
G(\lambda, Q)=D E_{\lambda}\left(\chi_{\lambda}+Q\right), \quad \text { where } \chi_{\lambda}=\left(\frac{2}{3}, q_{2, \lambda}, 0\right)
$$

By definition of the eigenvalue exchange solution, it holds

$$
\mathscr{G}(\lambda, 0)=0 .
$$

From Sect. 3, we know that

$$
D_{Q} \mathscr{G}(\lambda, 0)\left[h_{1}, h_{2}, h_{3}\right]=\left(\mathscr{M}_{\lambda}\left(h_{1}, h_{2}\right), \mathscr{L}_{\lambda} h_{3}\right),
$$

and $\mathscr{M}_{\lambda}$ is invertible since $v(\lambda)>0$. Recall indeed from Proposition 4.3 that the branch $\chi_{\lambda}$ is defined only when $v(\lambda)>0$.

As for $\mathscr{L}_{\lambda_{c}}$, its first eigenvalue is $\mu\left(\lambda_{c}\right)=0$ and it is simple. Therefore, we obtain

$$
\operatorname{dim} \operatorname{Ker} D_{Q} \mathscr{G}\left(\lambda_{c}, 0\right)=1=\operatorname{codim} \operatorname{Ran} D_{Q} \mathscr{G}\left(\lambda_{c}, 0\right)
$$

since $D_{Q} \mathscr{G}\left(\lambda_{c}, 0\right)$ is obviously Fredholm of index 0 .
Next we show that, for all $H \in \operatorname{Ker} D_{Q} \mathscr{G}\left(\lambda_{c}, 0\right), H \neq 0$, it holds

$$
\partial_{\lambda} D_{Q} \mathscr{G}\left(\lambda_{c}, 0\right) \cdot H \notin \operatorname{Ran} D_{Q} \mathscr{G}\left(\lambda_{c}, 0\right) .
$$

To this end, we use an argument similar to one in the proof of Alama et al. (2012, Theorem 5.24). Recall that

$$
\operatorname{Ker} D_{Q} \mathscr{G}\left(\lambda_{c}, 0\right)=\operatorname{Span}\left(0,0, h_{\lambda_{c}}\right),
$$

where $h_{\lambda}$ is an eigenfunction associated with the first eigenvalue of $\mathscr{L}_{\lambda}$ and can be chosen to depend smoothly on $\lambda$ (see the proof of Proposition 4.3).

Hence, it suffices to show that

$$
\left.\partial_{\lambda} \mathscr{L}_{\lambda}\right|_{\lambda=\lambda_{c}} h_{\lambda_{c}} \notin \operatorname{Ran} \mathscr{L}_{\lambda_{c}} .
$$

We obtain this latter fact as a consequence of $\mu^{\prime}\left(\lambda_{c}\right)<0$. Indeed, assume that there exists $h \in H_{0}^{1}$ such that

$$
\left.\partial_{\lambda} \mathscr{L}_{\lambda}\right|_{\lambda=\lambda_{c}} h_{\lambda_{c}}=\mathscr{L}_{\lambda_{c}} h .
$$

Then we compute, using the facts that $\mathscr{L}_{\lambda_{c}} h_{\lambda_{c}}=0$ and that $\mathscr{L}_{\lambda_{c}}$ is symmetric,

$$
\begin{aligned}
0>\mu^{\prime}\left(\lambda_{c}\right) & =\frac{d}{d \lambda}\left[<\mathscr{L}_{\lambda} h_{\lambda}, h_{\lambda}>\right]_{\lambda=\lambda_{c}} \\
& =<\left.\partial_{\lambda} \mathscr{L}_{\lambda}\right|_{\lambda=\lambda_{c}} h_{\lambda_{c}}, h_{\lambda_{c}}> \\
& =<\mathscr{L}_{\lambda_{c}} h, h_{\lambda_{c}}>=0,
\end{aligned}
$$

and we obtain a contradiction.
Thus, all the assumptions needed to apply Crandall-Rabinowitz' bifurcation theorem (Crandall and Rabinowitz 1971, Theorem 1.7) are satisfied: there exists a smooth function $\lambda(t)$ defined for small $t$, with $\lambda(0)=\lambda_{c}$, and a regular family $H_{t}=\left(h_{1, t}, h_{2, t}, h_{3, t}\right)$ taking values in $\left(0,0, h_{\lambda_{c}}\right)^{\perp} \subset H_{0}^{1}(-1,1)^{3}$ with $H_{0}=0$, such that, for any $Q$ close enough to $\chi_{\lambda_{c}}$,

$$
D E_{\lambda}(Q)=0 \Leftrightarrow\left\{\begin{array}{c}
\quad Q=\chi_{\lambda} \\
\text { or } \quad \lambda=\lambda(t) \text { and } Q=\chi_{\lambda(t)}+t\left(0,0, h_{\lambda_{c}}\right)+t^{2} H_{t}
\end{array}\right.
$$

One can say a little bit more about the new branch of solutions thus obtained. Indeed, changing $q_{3}$ to $-q_{3}$ leaves the equations (6) invariant. More precisely, given $Q=\left(q_{1}, q_{2}, q_{3}\right)$ a solution of (6), the map $\widetilde{Q}=\left(q_{1}, q_{2},-q_{3}\right)$ is automatically a solution of (6). In particular, to a solution

$$
Q=\chi_{\lambda(t)}+t\left(0,0, h_{\lambda_{c}}\right)+t^{2} H_{t},
$$

corresponds a solution

$$
\widetilde{Q}=\chi_{\lambda(t)}-t\left(0,0, h_{\lambda_{c}}\right)+t^{2} \widetilde{H}_{t} .
$$

Since both $Q$ and $\widetilde{Q}$ are close to $\chi_{\lambda_{c}}$, we deduce that

$$
\lambda(t)=\lambda(-t) \quad \text { and } h_{1,-t}=h_{1, t}, h_{2,-t}=h_{2, t}, h_{3,-t}=-h_{3, t} .
$$

This ends the proof of Theorem 5.2.

## 6 The Perturbed Case $\theta \approx-8$

Now we turn back to the case of a general temperature $\theta \in(-\infty, 1]$. A closer look at the proof in subsection 5.2 will convince us that a result similar to Theorem 5.2 holds for any $\theta$ satisfying some non-degeneracy assumptions. After having checked that these non-degeneracy assumptions are stable under small perturbations of $\theta$, we will obtain as a corollary a result similar to Theorem 5.2 in the perturbed case $\theta \approx-8$.

Theorem 6.1 Assume that $\theta$ is such that the branch of eigenvalue exchange solutions $\lambda \mapsto \chi_{\lambda}$ given by Proposition 4.3 has the following two properties:
(i) there exists $\lambda \in\left(0, \lambda_{*}\right)$ such that $\mu(\lambda)<0$.
(ii) denoting by $\lambda_{c}>0$ the infimum of all such $\lambda$ :

$$
\lambda_{c}=\inf \left\{\lambda \in\left(0, \lambda_{*}\right): \mu(\lambda)<0\right\}
$$

it holds

$$
\mu^{\prime}\left(\lambda_{c}\right)<0 .
$$

Then, there exist $\delta, \varepsilon>0$ and a neighborhood $\mathfrak{A}$ of $\chi_{\lambda_{c}}$ in $\mathscr{H}$, such that the solutions of

$$
D E_{\lambda}(Q)=0, \quad(\lambda, Q) \in\left(\lambda_{c}-\delta, \lambda_{c}+\delta\right) \times \mathfrak{A}
$$

are exactly

$$
Q=\chi_{\lambda} \quad \text { or } \quad\left\{\begin{array}{l}
\lambda=\lambda(t) \\
Q=\chi_{\lambda_{c}}+t\left(0,0, h_{c}\right)+t^{2} H_{t},
\end{array} \quad \text { for some } t \in(-\varepsilon, \varepsilon),\right.
$$

where $\lambda(t) \in\left(\lambda_{c}-\delta, \lambda_{c}+\delta\right)$ and $H_{t} \in h_{c}^{\perp}$ are smooth functions of $t \in(-\varepsilon, \varepsilon)$. Moreover, the following symmetry properties are satisfied:

$$
\begin{equation*}
\lambda(-t)=\lambda(t), \quad \text { and } \quad h_{1,-t}=h_{1, t}, h_{2,-t}=h_{2, t}, h_{3,-t}=-h_{3, t}, \tag{37}
\end{equation*}
$$

where $H_{t}$ is identified with $\left(h_{1, t}, h_{2, t}, h_{3, t}\right)$ via (1).
Proof Looking at the proof of Theorem 5.2 in subsection 5.2, we see that we have really only used the facts that for $\theta=-8,(i)$ and (ii) are satisfied. Hence, the proof of Theorem 5.2 may be reproduced word for word to prove Theorem 6.1.

Theorem 6.1 is an abstract theorem: if $\theta$ satisfies some conditions, then we have a concrete result. But it does not tell us anything about the validity of such conditions in general.

Let us say a few words about these conditions. In view of Proposition 4.3, $\lambda_{*}$ can be interpreted as the point where the eigenvalue exchange solution looses its stability with respect to symmetry-preserving perturbations. Condition (i) asks for $\mu(\lambda)$ to become negative before this point is reached. That is, condition (i) could be rephrased as: as $\lambda$ grows, starting from the unique solution for small $\lambda$, a symmetry-breaking loss of stability occurs before a possible symmetry-preserving loss of stability. And condition (ii) asks for the symmetry-breaking loss of stability to be non-degenerate. Hence, condition (ii) is typically a generic condition.

Remark that in the special case $\theta=-8$, we have shown (Theorem 5.1) that symmetry-preserving loss of stability does not occur at all, and that symmetry-breaking loss of stability does occur, in a non-degenerate way. Now we are going to show that these conditions propagate to nearby $\theta$. This is the content of the next result.

Proposition 6.2 If $\theta_{0}<1$ satisfies conditions (i) and (ii) of Theorem 6.1, then there exists $\varepsilon>0$ such that every $\theta<1$ with $\left|\theta-\theta_{0}\right|<\varepsilon$ also satisfies (i) and (ii).

Proof During this proof we will emphasize the dependence on $\theta$ of the objects we have been working with. For instance we will write $\mathscr{H}_{\theta}, E_{\theta, \lambda}^{e e}, \lambda_{*}(\theta)$, and so on.

Let us start by remarking that a value $\lambda_{0}$ (provided by Proposition 4.2), under which there is uniqueness of the solution, may be chosen independently of $\theta$ in a neighborhood of $\theta_{0}$. Indeed, the proof of Lemma 8.1 shows that this value of $\lambda_{0}$ depends on the $W^{2, \infty}$ norm of the bulk energy density $f$ restricted to values of $Q$ satisfying the maximum principle (5). It is clear from the expression of $f$ and (5) that this $W^{2, \infty}$ norm depends at least continuously on $\theta$. We may thus choose a $\lambda_{0}$ that works for all $\theta$ in a fixed neighborhood of $\theta_{0}$.

The idea of the proof is to use the implicit function theorem to define eigenvalue exchange solutions $\chi_{\lambda, \theta}$ depending smoothly on $\lambda$ and $\theta$. For $\theta$ close enough to $\theta_{0}$, this branch will look very much like the branch $\chi_{\lambda, \theta_{0}}$, and thus will satisfy $(i)$ and (ii).

To apply the implicit function theorem we need a fixed space, but $\mathscr{H}_{\theta}^{\text {ee }}$ depends on $\theta$. Thus, we fix $\chi_{\theta} \in \mathscr{H}_{\theta}^{\text {ee }}$ depending smoothly on $\theta$ (for instance take $\chi_{\theta}$ to be affine), and we will work instead in the space $H_{0}^{1}(-1,1)^{2}$ after having translated by $\chi_{\theta}$.

Since we will apply the implicit function theorem near each $\lambda$, but need to obtain for each $\theta$ a whole branch $\lambda \mapsto \chi_{\lambda, \theta}$, we will have to restrict $\lambda$ to a compact interval. That is why we choose $\lambda_{1} \in\left(\lambda_{c}, \lambda_{*}\right)$, where $\lambda_{c}=\lambda_{c}\left(\theta_{0}\right)($ defined by $(i i))$ and $\lambda_{*}=\lambda_{*}\left(\theta_{0}\right)$.

We consider the smooth function $\mathscr{F}$ defined by

$$
\mathscr{F}(\theta, \lambda, \chi)=D E_{\lambda, \theta}^{\mathrm{ee}}\left(\chi_{\theta}+\chi\right) \in H^{-1}(-1,1)^{2}
$$

for $\theta<1, \lambda>0$ and $\chi \in H_{0}^{1}(-1,1)^{2}$. For all $\lambda \in\left(0, \lambda_{*}\right)$, it holds

$$
\mathscr{F}\left(\theta_{0}, \lambda, \chi_{\lambda, \theta_{0}}-\chi_{\theta_{0}}\right)=0,
$$

and, since $v(\lambda)>0$, the partial differential

$$
D_{\chi} \mathscr{F}\left(\theta_{0}, \lambda, \chi_{\lambda, \theta_{0}}-\chi_{\theta_{0}}\right) \text { is invertible. }
$$

Hence the implicit function theorems provide us with $\varepsilon_{\lambda}>0$ and $\mathfrak{A}_{\lambda}$ a neighborhood of $\chi_{\lambda, \theta_{0}}-\chi_{\theta_{0}}$ such that the equation

$$
\mathscr{F}\left(\theta, \lambda^{\prime}, \chi\right)=0, \quad\left|\theta-\theta_{0}\right|<\varepsilon_{\lambda}, \quad\left|\lambda^{\prime}-\lambda\right|<\varepsilon_{\lambda}, \quad \chi \in \mathfrak{A}_{\lambda},
$$

has a unique solution $\chi_{\lambda^{\prime}, \theta}-\chi_{\theta}$ depending smoothly on $\left(\lambda^{\prime}, \theta\right)$.
Using the compactness of $\left[\lambda_{0} / 2, \lambda_{1}\right]$, we deduce the existence of $\varepsilon>0$ and $\mathfrak{A}$ a neighborhood of 0 in $H_{0}^{1}(-1,1)^{2}$, such that the equation

$$
\mathscr{F}(\theta, \lambda, \chi)=0, \quad\left|\theta-\theta_{0}\right|<\varepsilon, \frac{\lambda_{0}}{2} \leq \lambda \leq \lambda_{1}, \chi \in \chi_{\lambda, \theta_{0}}-\chi_{\theta_{0}}+\mathfrak{A},
$$

has a unique solution $\chi_{\theta, \lambda}-\chi_{\theta}$ which depends smoothly on $(\theta, \lambda)$. Hence, for every $\theta \in\left(\theta_{0}-\varepsilon, \theta_{0}+\varepsilon\right)$, the unique smooth branch of eigenvalue exchange solutions given by Proposition 4.3 is $\lambda \mapsto \chi_{\theta, \lambda}$, defined at least up to $\lambda_{1}$, and it depends smoothly on $\theta$.

More precisely, we have just proven that $\chi_{\theta, \lambda}$ depends smoothly on $(\theta, \lambda) \in\left(\theta_{0}-\right.$ $\left.\varepsilon, \theta_{0}+\varepsilon\right) \times\left(\lambda_{0} / 2, \lambda_{1}\right)$. On the other hand, since $\lambda_{0}$ is chosen in such a way that the unique solution $\chi_{\theta, \lambda}$ is non-degenerate for $\lambda<\lambda_{0}$ (see the proof of Theorem 8.1), we may apply the implicit function theorem to obtain that $(\theta, \lambda) \mapsto \chi_{\theta, \lambda}$ is smooth also for small $\lambda$. Hence $\chi_{\theta, \lambda}$ depends smoothly on $(\theta, \lambda) \in\left(\theta_{0}-\varepsilon, \theta_{0}+\varepsilon\right) \times\left(0, \lambda_{1}\right)$.

Recall that, for fixed $\theta$, given a branch of EE solutions $\chi_{\lambda}$, we have defined $\mu(\lambda)$ in (13), as the first eigenvalue of the free energy second variation around $\chi_{\lambda}$ with respect to symmetry-breaking perturbations. Here, we emphasize the dependence on $\theta$ by writing $\mu(\theta, \lambda)$. That is, $\mu(\theta, \lambda)$ is the first eigenvalue of $\mathscr{L}_{\theta, \lambda}$, which is the linear operator associated to the quadratic form $D^{2} E_{\theta, \lambda}$ restricted to the space $H_{\mathrm{sb}}$ of symmetry-breaking perturbations (see Sect. 3).

Since $(\theta, \lambda) \mapsto \chi_{\theta, \lambda}$ is smooth, we prove, exactly as in Proposition 4.3 for $\lambda \mapsto$ $\mu(\lambda)$, that

$$
\left(\theta_{0}-\varepsilon, \theta_{0}+\varepsilon\right) \times\left(0, \lambda_{1}\right) \ni(\theta, \lambda) \mapsto \mu(\theta, \lambda)
$$

is smooth.
In particular, since—by $(i)$-there exists $\lambda_{2} \in\left(\lambda_{c}\left(\theta_{0}\right), \lambda_{1}\right)$ such that $\mu\left(\theta_{0}, \lambda_{2}\right)<0$, it follows that we may chose $\varepsilon$ small enough, so that

$$
\mu\left(\theta, \lambda_{2}\right)<0 \quad \forall \theta \in\left(\theta_{0}-\varepsilon, \theta_{0}+\varepsilon\right)
$$

i.e., $(i)$ is satisfied for $\theta$ close enough to $\theta_{0}$.

By definition of $\lambda_{c}=\lambda_{c}\left(\theta_{0}\right)$, and since $\theta_{0}$ satisfies (ii), it holds

$$
\mu\left(\theta_{0}, \lambda_{c}\right)=0, \quad \frac{\partial \mu}{\partial \lambda}\left(\theta_{0}, \lambda_{c}\right)>0 .
$$

Therefore the implicit function theorem ensures the existence of a smooth map $\lambda(\theta)$ defined—up to choosing $\varepsilon$ small enough—for $\theta \in\left(\theta_{0}-\varepsilon, \theta_{0}+\varepsilon\right)$, such that

$$
\mu\left(\theta_{0}, \lambda(\theta)\right)=0, \quad \frac{\partial \mu}{\partial \lambda}(\theta, \lambda(\theta))>0 .
$$

In order to complete the proof of Proposition 6.2, we need to show that this $\lambda(\theta)$ is really the critical value $\lambda_{c}(\theta)$ that appears in (ii).

That is, we need to prove that (for $\varepsilon$ small enough),

$$
\begin{equation*}
\mu(\theta, \lambda)>0 \text { for } \theta \in\left(\theta_{0}-\varepsilon, \theta_{0}+\varepsilon\right), \lambda \in(0, \lambda(\theta)) . \tag{38}
\end{equation*}
$$

We start by noting that the choice of $\lambda_{0}$ in the proof of the uniqueness result Theorem 8.1 can be such that

$$
\begin{equation*}
\mu(\theta, \lambda) \geq c_{0} \quad \forall(\theta, \lambda) \in\left(\theta_{0}-\varepsilon, \theta_{0}+\varepsilon\right) \times\left(0, \lambda_{0}\right), \tag{39}
\end{equation*}
$$

for some $c_{0}>0$. On the other hand, $\varepsilon$ may be chosen in such a way that it holds

$$
\begin{equation*}
\frac{\partial \mu}{\partial \lambda}(\theta, \lambda)>0 \quad \text { for }(\theta, \lambda) \in\left(\theta_{0}-\varepsilon, \theta_{0}+\varepsilon\right) \times\left(\lambda_{c}-\delta, \lambda_{c}+\delta\right) \tag{40}
\end{equation*}
$$

Using the compactness of $\left[\lambda_{0}, \lambda_{c}-\delta\right]$ and the fact that $\mu\left(\theta_{0}, \lambda\right)>0$ for all $\lambda \in\left(0, \lambda_{c}\right)$, we may also choose $\varepsilon$ such that we have

$$
\begin{equation*}
\mu(\theta, \lambda)>0 \quad \forall(\theta, \lambda) \in\left(\theta_{0}-\varepsilon, \theta_{0}+\varepsilon\right) \times\left[\lambda_{0}, \lambda_{c}-\delta\right] . \tag{41}
\end{equation*}
$$

Putting together (39), (6) and (41), we obtain (38). Therefore, $\lambda(\theta)$ is really the infimum of those $\lambda$ for which $\mu(\theta, \lambda)<0$, and $\theta \in\left(\theta_{0}-\varepsilon, \theta_{0}+\varepsilon\right)$ satisfies (ii).

As we pointed out at the beginning of the present section, a corollary of Theorem 5.1 and Proposition 6.2 is that the bifurcation result Theorem 6.1 applies to all $\theta$ close enough to the special value $\theta=-8$.

Corollary 6.3 There exists $\varepsilon>0$ such that, for any $\theta<1$ with $|\theta+8|<\varepsilon$, a symmetric pitchfork bifurcation occurs from the branch of eigenvalue exchange solutions starting at small $\lambda$, in the sense that Theorem 6.1 applies.

## Appendix 1: Principle of Symmetric Criticality

Proposition 7.1 Let H be a Hilbert space, G a group acting linearly and isometrically on $H$ and $\Sigma=H^{G}$ the subspace of symmetric elements (that is, $x \in \Sigma$ iff $g x=x$ $\forall g \in G)$. Let $f: H \rightarrow \mathbb{R}$ be a $G$-invariant $C^{1}$ function. It holds:
(i) If $x \in \Sigma$ is a critical point of $f_{\mid \Sigma}$, then $x$ is a critical point of $f$.
(ii) If in addition $f$ is $C^{2}$, it further holds

$$
D^{2} f(x) \cdot h \cdot k=0 \text { for } h \in \Sigma, k \in \Sigma^{\perp}
$$

i.e., the orthogonal decomposition $H=\Sigma \oplus \Sigma^{\perp}$ is also orthogonal for the bilinear form $D^{2} f(x)$.

Item (i) of the above proposition is only a particularly simple case of Palais' Principle of symmetric criticality (Palais 1979). Item (ii), however, does not seem to be explicitly stated in the literature-as far as we know. Using the same tools as in Sect. 2 of Palais (1979), it is not hard to see that an equivalent of (ii) is actually valid if $H$ is replaced by a Riemannian manifold $\mathscr{M}$ on which the group $G$ acts isometrically. In this case, $\Sigma$ is a submanifold of $\mathscr{M}$ and, at a symmetric critical point $x$, the orthogonal decomposition

$$
T_{x} \mathscr{M}=T_{x} \Sigma \oplus\left(T_{x} \Sigma\right)^{\perp}
$$

is also orthogonal for the bilinear form $D^{2} f(x)$.
Proof of Proposition 7.1 As already pointed out, item (i) is a particular case of Palais (1979, Sect. 2). We nevertheless present a complete Proof of Proposition 7.1 here, since in the simple framework we consider, the proof of $(i)$ is really straightforward.

The fact that $f$ is $G$-invariant means that it holds

$$
\begin{equation*}
f(g x)=f(x) \quad \forall g \in G, x \in H . \tag{42}
\end{equation*}
$$

Since the action of $G$ on $H$ is linear, differentiating (42) we obtain

$$
\begin{equation*}
D f(g x) \cdot g h=D f(x) \cdot h \quad \forall h \in H . \tag{43}
\end{equation*}
$$

Applying (43) for a symmetric $x$, i.e., $x \in \Sigma$, we have

$$
<\nabla f(x), g h>=<\nabla f(x), h>\quad \forall h \in H .
$$

Note that here we distinguish between the differential $D f(x) \in H^{*}$ and the gradient $\nabla f(x) \in H$. Similarly, below we will distinguish between the second-order differential $D^{2} f(x) \in \mathscr{L}\left(H, H^{*}\right)$ and the Hessian $\nabla^{2} f(x) \in \mathscr{L}(H)$. Since $g$ is a linear isometry, we conclude that

$$
\begin{equation*}
g^{-1} \nabla f(x)=\nabla f(x) \quad \forall g \in G, \quad \text { i.e., } \nabla f(x) \in \Sigma . \tag{44}
\end{equation*}
$$

Therefore, if we know in addition that $x$ is a critical point of $f_{\mid \Sigma}$, which means that $\nabla f(x) \in \Sigma^{\perp}$, it must hold $\nabla f(x)=0$. This proves $(i)$.

Now assume that $f$ is $C^{2}$ and differentiate (43) to obtain

$$
\begin{equation*}
D^{2} f(g x) \cdot g h \cdot g k=D^{2} f(x) \cdot h \cdot k \quad \forall h, k \in H . \tag{45}
\end{equation*}
$$

In particular, if $x$ and $h$ are symmetric (i.e., belong to $\Sigma$ ), and if we denote by $\nabla^{2} f(x)$ the Hessian of $f$ at $x$, (45) becomes

$$
<g^{-1}\left(\nabla^{2} f(x) h\right), k>=<\nabla^{2} f(x) h, k>\quad \forall x \in \Sigma, h \in \Sigma, k \in H,
$$

so that $\nabla^{2} f(x) h$ is symmetric. Hence it is orthogonal to any $k \in \Sigma^{\perp}$, which proves (ii).

## Appendix 2: Uniqueness of Critical Points for Small $\lambda$

Let $\Omega \subset \mathbb{R}^{N}$ be a smooth bounded domain, and $f: \mathbb{R}^{d} \rightarrow \mathbb{R}$ a $W_{\text {loc }}^{2, \infty}$ map. We are interested in critical points of functionals of the form

$$
\begin{equation*}
E_{\lambda}(u)=\int_{\Omega} \frac{1}{2 \lambda^{2}}|\nabla u|^{2}+\int_{\Omega} f(u), \tag{46}
\end{equation*}
$$

i.e., solutions $u \in H^{1}(\Omega)^{d}$ of the equation

$$
\begin{equation*}
\Delta u=\lambda^{2} \nabla f(u) \quad \text { in } \mathscr{D}^{\prime}(\Omega) . \tag{47}
\end{equation*}
$$

Note that (47) implies in particular that $\nabla f(u) \in L_{l o c}^{1}$.
We prove the following:

Theorem 8.1 Assume that there exists $C>0$ such that $\nabla f(x) \cdot x \geq 0$ for any $x \in \mathbb{R}^{d}$ with $|x| \geq C$.

Let $g \in L^{\infty} \cap H^{1 / 2}(\partial \Omega)^{d}$. There exists $\lambda_{0}=\lambda_{0}(\Omega, f, g)$ such that, for any $\lambda \in\left(0, \lambda_{0}\right), E_{\lambda}$ admits at most one critical point with $\operatorname{tr} u=g$ on $\partial \Omega$.

Theorem 8.1 is a direct consequence of Lemmas 8.2 and 8.3 below. Indeed, Lemma 8.2 ensures that, for sufficiently small $\lambda, E_{\lambda}$ admits at most one critical point satisfying a given $L^{\infty}$ bound (independent of $\lambda$ ). And in Lemma 8.3 we prove that the assumption on $f$ implies such a bound for critical points of $E_{\lambda}$.

Lemma 8.2 Let $C>0$. There exists $\lambda_{0}=\lambda_{0}(C, f, \Omega)$ such that, for any $\lambda \in\left(0, \lambda_{0}\right)$ and any $g \in H^{1 / 2}(\partial \Omega)^{d}, E_{\lambda}$ admits at most one critical point $u$ satisfying $|u| \leq C$ a.e. and $\operatorname{tr} u=g$.

Proof Let

$$
X:=\left\{u \in H^{1}(\Omega):|u| \leq C \text { a.e. }\right\} .
$$

We show that, for $\lambda$ small enough, $E_{\lambda}$ is strictly convex on $X$.

Let $u, v \in X$. Then $u-v \in H_{0}^{1}(\Omega)^{d}$. Using Poincaré's inequality, we obtain

$$
\begin{align*}
E_{\lambda}\left(\frac{u+v}{2}\right)= & \frac{1}{8 \lambda^{2}} \int|\nabla u+\nabla v|^{2}+\int f\left(\frac{u+v}{2}\right) \\
= & \frac{1}{4 \lambda^{2}} \int|\nabla u|^{2}+\frac{1}{4 \lambda^{2}} \int|\nabla v|^{2}-\frac{1}{8 \lambda^{2}} \int|\nabla(u-v)|^{2}+\int f\left(\frac{u+v}{2}\right) \\
= & \frac{1}{2} E_{\lambda}(u)+\frac{1}{2} E_{\lambda}(v)-\frac{1}{8 \lambda^{2}} \int|\nabla(u-v)|^{2} \\
& +\int\left[f\left(\frac{u+v}{2}\right)-\frac{1}{2} f(u)-\frac{1}{2} f(v)\right] \\
\leq & \frac{1}{2} E_{\lambda}(u)+\frac{1}{2} E_{\lambda}(v)-\frac{c_{1}(\Omega)}{\lambda^{2}}\|u-v\|_{L^{2}}^{2} \\
& +\int\left[f\left(\frac{u+v}{2}\right)-\frac{1}{2} f(u)-\frac{1}{2} f(v)\right] . \tag{48}
\end{align*}
$$

On the other hand, for any $x, y \in \mathbb{R}^{d}$ satisfying $|x|,|y| \leq C$, it holds

$$
\begin{equation*}
f\left(\frac{x+y}{2}\right)-\frac{1}{2} f(x)-\frac{1}{2} f(y) \leq\|f\|_{W^{2, \infty}\left(B_{C}\right)}|x-y|^{2} . \tag{49}
\end{equation*}
$$

Plugging (49) into (48) we obtain, for some $c_{2}=c_{2}(\Omega, f, C)>0$,

$$
\begin{aligned}
E_{\lambda}\left(\frac{u+v}{2}\right) & \leq \frac{1}{2} E_{\lambda}(u)+\frac{1}{2} E_{\lambda}(v)-\frac{c_{1}}{\lambda^{2}}\|u-v\|_{L^{2}}^{2}+c_{2}\|u-v\|_{L^{2}}^{2} \\
& =\frac{1}{2} E_{\lambda}(u)+\frac{1}{2} E_{\lambda}(v)-\frac{c_{1}}{2 \lambda^{2}}\|u-v\|_{L^{2}}^{2}-c_{2}\left(\frac{c_{1}}{2 c_{2} \lambda^{2}}-1\right)\|u-v\|_{L^{2}}^{2} .
\end{aligned}
$$

Hence, for $\lambda \leq \lambda_{0}:=\sqrt{c_{1} /\left(2 c_{2}\right)}$, it holds

$$
E_{\lambda}\left(\frac{u+v}{2}\right)<\frac{1}{2} E_{\lambda}(u)+\frac{1}{2} E_{\lambda}(v) \quad \forall u, v \in X, u \neq v .
$$

Thus, $E_{\lambda}$ is strictly convex on $X$.
To conclude the proof, assume that for a $\lambda \in\left(0, \lambda_{0}\right)$, there exist two solutions $u_{1}$ and $u_{2}$ of (47), belonging to $X$. Then one easily shows that [0, 1] $\ni t \mapsto E_{\lambda}\left(t u_{1}+(1-t) u_{2}\right)$ is $C^{1}$ and that its derivative vanishes at 0 and 1 , which is incompatible with the strict convexity of $E_{\lambda}$.

Lemma 8.3 Assume that there exists $C>0$ such that

$$
|x| \geq C \Rightarrow \nabla f(x) \cdot x \geq 0
$$

Let $g \in L^{\infty} \cap H^{1 / 2}(\partial \Omega)^{d}$. If $u \in H_{g}^{1}(\Omega)^{d}$ is a critical point of $E_{\lambda}$, then it holds

$$
|u| \leq \max \left(C,\|g\|_{\infty}\right) \quad \text { a.e. }
$$

Proof We may assume $C=\max \left(C,\|g\|_{\infty}\right)>0$.
Let $\varphi \in C^{\infty}(\mathbb{R})$ be such that:

$$
\left\{\begin{array}{l}
\varphi \geq 0,  \tag{50}\\
\varphi^{\prime} \geq 0, \\
\varphi(t)=0 \quad \text { for } t \leq C^{2}, \\
\varphi(t)=1 \quad \text { for } t \geq T, \text { for some } T>C^{2}
\end{array}\right.
$$

Let $w=\varphi\left(|u|^{2}\right)$. The assumptions on $\varphi$ ensure that $w \geq 0$, and $w=0$ in $\{|u| \leq C\}$. Therefore, taking the scalar product of (47) with $w u$ and using the assumption that $\nabla f(u) \cdot u \geq 0$ outside of $\{|u| \leq C\}$, we obtain

$$
\begin{equation*}
\frac{1}{\lambda^{2}} w u \cdot \Delta u=w \nabla f(u) \cdot u \geq 0 \quad \text { a.e. } \tag{51}
\end{equation*}
$$

Since $w u \in H_{0}^{1}(\Omega)^{d}$, we may apply Lemma 8.4 below, to deduce

$$
\begin{equation*}
\int_{\Omega} \nabla u \cdot \nabla(w u) \leq 0 . \tag{52}
\end{equation*}
$$

On the other hand, it holds

$$
\int_{\Omega} \nabla u \cdot \nabla(w u)=\int_{\Omega} w|\nabla u|^{2}+\int_{\Omega} 2 \sum_{k}\left(u \cdot \partial_{k} u\right)^{2} \varphi^{\prime}\left(|u|^{2}\right),
$$

so that we have in fact

$$
\begin{equation*}
\int_{\Omega} w|\nabla u|^{2} \leq 0 . \tag{53}
\end{equation*}
$$

Finally we may choose an increasing sequence $\varphi_{k}$ of smooth maps satisfying (50) and converging to $\mathbf{1}_{t>C^{2}}$. Then, $w_{k}=\varphi_{k}\left(|u|^{2}\right)$ is increasing and converges a.e. to $\mathbf{1}_{|u|>C}$, and we conclude that

$$
\int_{|u|>C}|\nabla u|^{2}=0,
$$

so that $|u| \leq C$ a.e.
The following result, which we used in the proof of Lemma 8.3, is due to Pierre Bousquet.

Lemma 8.4 Let $u \in H^{1}(\Omega)^{d}$ and assume that $\Delta u=g \in L_{l o c}^{1}(\Omega)^{d}$. Then, for any $\zeta \in H_{0}^{1}(\Omega)^{d}$,

$$
\begin{equation*}
\zeta \cdot g \geq 0 \text { a.e. } \Longrightarrow \int \nabla \zeta \cdot \nabla u \leq 0 . \tag{54}
\end{equation*}
$$

Proof We proceed in three steps: first we show that (54) is valid for $\zeta \in H^{1} \cap L^{\infty}(\Omega)^{d}$ with compact support in $\Omega$, then for $\zeta \in H_{0}^{1} \cap L^{\infty}(\Omega)^{d}$, and eventually for $\zeta \in$ $H_{0}^{1}(\Omega)^{d}$.

Step 1: $\zeta \in H_{c}^{1} \cap L^{\infty}$.
Since $\zeta$ is bounded and compactly supported, there exists a sequence $\zeta_{k}$ of $C_{c}^{\infty}$ functions, a constant $C>0$, and a compact $K \subset \Omega$, such that

$$
\operatorname{supp} \zeta_{k} \subset K, \quad\left\|\zeta_{k}\right\|_{\infty} \leq C, \text { and } \zeta_{k} \longrightarrow \zeta \text { in } H^{1} \text { and a.e. }
$$

Since $\zeta_{k} \in C_{c}^{\infty}(\Omega)^{d}$, it holds, by definition of the weak laplacian,

$$
\int \zeta_{k} \cdot g=-\int \nabla \zeta_{k} \cdot \nabla u
$$

and we may pass to the limit (using dominated convergence on the compact $K$ for the left-hand side) to obtain

$$
\int_{\Omega} \zeta \cdot g=-\int \nabla \zeta \cdot \nabla u
$$

which implies (54).
Step 2: $\zeta \in H_{0}^{1} \cap L^{\infty}$.
Let $\theta_{k} \in C_{c}^{\infty}(\Omega)$ be such that

$$
0 \leq \theta_{k} \leq 1, \theta_{k}(x)=1 \text { if } d(x, \partial \Omega)>\frac{1}{k}, \text { and }\left|\nabla \theta_{k}(x)\right| \leq \frac{c}{d(x, \partial \Omega)}
$$

and define $\zeta_{k}=\theta_{k} \zeta \in H_{c}^{1} \cap L^{\infty}(\Omega)^{d}$.
Assuming that $\zeta \cdot g \geq 0$ a.e., we deduce that $\zeta_{k} \cdot g \geq 0$ a.e., and thus we may apply Step 1 to $\zeta_{k}$ : it holds

$$
\begin{equation*}
0 \geq \int \nabla \zeta_{k} \cdot \nabla u=\int \theta_{k} \nabla \zeta \cdot \nabla u+\int \nabla \theta_{k} \cdot \nabla u \cdot \zeta . \tag{55}
\end{equation*}
$$

The first term in the right-hand side of (55) converges to $\int \nabla \zeta \cdot \nabla u$, by dominated convergence. Therefore we only need to prove that the second term in the right-hand side of (55) converges to zero. To this end we use the following Hardy-type inequality:

$$
\begin{equation*}
\int \frac{|\zeta|^{2}}{d(x, \partial \Omega)^{2}} \leq C \int|\nabla \zeta|^{2}, \quad \forall \zeta \in H_{0}^{1}(\Omega) \tag{56}
\end{equation*}
$$

Using (56) and the Hölder inequality, we obtain

$$
\left|\int \nabla \theta_{k} \cdot \nabla u \cdot \zeta\right|^{2} \leq C\|\nabla \zeta\|_{L^{2}}^{2} \int_{d(x, \partial \Omega)>1 / k}|\nabla u|^{2} \longrightarrow 0
$$

which concludes the proof of Step 2.
Step 3: $\zeta \in H_{0}^{1}$.

We define $\zeta_{k}=P_{k}(\zeta)$, where $P_{k}: \mathbb{R}^{d} \rightarrow \mathbb{R}^{d}$ is given by

$$
P_{k}(x)= \begin{cases}x & \text { if }|x| \leq k \\ \frac{k}{|x|} x & \text { if }|x|>k\end{cases}
$$

Then $\zeta_{k} \in H_{0}^{1} \cap L^{\infty}(\Omega)$ and $\zeta_{k} \rightarrow \zeta$ in $H^{1}$.
If $\zeta \cdot g \geq 0$, then it obviously holds $\zeta_{k} \cdot g \geq 0$, so that we may apply Step 2 to $\zeta_{k}$ and obtain

$$
\int \nabla \zeta_{k} \cdot \nabla u \leq 0
$$

Letting $k$ go to $\infty$ in this last inequality provides the desired conclusion.

## Appendix 3: Second Variation of the Energy

At a map $Q \in H^{1}(-1,1)^{3}$, the second variation of the energy reads

$$
D^{2} E(Q)[H]=\int\left(\frac{1}{\lambda^{2}}\left(H^{\prime}\right)^{2}+D^{2} f(Q)[H]\right) \mathrm{d} x
$$

where

$$
D^{2} f(Q)[H]=\frac{\theta}{3}|H|^{2}-4 Q \cdot H^{2}+(Q \cdot H)^{2}+\frac{1}{2}|Q|^{2}|H|^{2} .
$$

If we take $Q=\chi=\left(q_{1}, q_{2}, 0\right)$, and consider separately perturbations $H_{s p}=$ $\left(h_{1}, h_{2}, 0\right)$ and $H_{\mathrm{sb}}=\left(0,0, h_{3}\right)$, we have

$$
\begin{aligned}
\left|H_{\mathrm{sp}}\right|^{2} & =6 h_{1}^{2}+2 h_{2}^{2} & \left|H_{\mathrm{sb}}\right|^{2} & =2 h_{3}^{2} \\
\chi \cdot H_{\mathrm{sp}}^{2} & =2 q_{1}\left(h_{2}^{2}-3 h_{1}^{2}\right)+4 q_{2} h_{1} h_{2} & \chi \cdot H_{\mathrm{sb}}^{2} & =2 q_{1} h_{3}^{2} \\
\chi \cdot H_{\mathrm{sp}} & =6 q_{1} h_{1}+2 q_{2} h_{2} & \chi \cdot H_{\mathrm{sb}} & =0,
\end{aligned}
$$

so that we can compute

$$
\begin{aligned}
D^{2} f(\chi)\left[H_{\mathrm{sp}}\right]= & 6\left(\frac{\theta}{3}+2 q_{1}+9 q_{1}^{2}+q_{2}^{2}\right) h_{1}^{2} \\
& +2\left(\frac{\theta}{3}-4 q_{1}+3 q_{1}^{2}+3 q_{2}^{2}\right) h_{2}^{2} \\
& +8 q_{2}\left(3 q_{1}-2\right) h_{1} h_{2} \\
D^{2} f(\chi)\left[H_{\mathrm{sb}}\right]= & 2\left(\frac{\theta}{3}-4 q_{1}+3 q_{1}^{2}+q_{2}^{2}\right) h_{3}^{2} .
\end{aligned}
$$

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