# SMOOTH APPROXIMATIONS WITHOUT CRITICAL POINTS OF CONTINUOUS MAPPINGS BETWEEN BANACH SPACES, AND DIFFEOMORPHIC EXTRACTIONS OF SETS 

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#### Abstract

Let $E, F$ be separable Hilbert spaces, and assume that $E$ is infinite-dimensional. We show that for every continuous mapping $f: E \rightarrow F$ and every continuous function $\varepsilon: E \rightarrow(0, \infty)$ there exists a $C^{\infty}$ mapping $g: E \rightarrow F$ such that $\|f(x)-g(x)\| \leq \varepsilon(x)$ and $D g(x): E \rightarrow F$ is a surjective linear operator for every $x \in E$. We also provide a version of this result where $E$ can be replaced with a Banach space from a large class (including all the classical spaces with smooth norms, such as $c_{0}, \ell_{p}$ or $\left.L^{p}, 1<p<\infty\right)$, and $F$ can be taken to be any Banach space such that there exists a bounded linear operator from $E$ onto $F$. In particular, for such $E, F$, every continuous mapping $f: E \rightarrow F$ can be uniformly approximated by smooth open mappings. Part of the proof provides results of independent interest that improve some known theorems about diffeomorphic extractions of closed sets from Banach spaces or Hilbert manifolds.


## 1. Introduction and main results

The main purpose of this paper is to show the following two results.
Theorem 1.1. Let E, F be separable Hilbert spaces, and assume that $E$ is infinite-dimensional. Then, for every continuous mapping $f: E \rightarrow F$ and every continuous function $\varepsilon: E \rightarrow(0, \infty)$ there exists a $C^{\infty}$ mapping $g: E \rightarrow F$ such that $\|f(x)-g(x)\| \leq \varepsilon(x)$ and $D g(x): E \rightarrow F$ is a surjective linear operator for every $x \in E$.

Theorem 1.2. Let $E$ be one of the classical Banach spaces $c_{0}$, $\ell_{p}$ or $L^{p}, 1<p<\infty$. Let $F$ be a Banach space, and assume that there exists a bounded linear operator from $E$ onto $F$. Then, for every continuous mapping $f: E \rightarrow F$ and every continuous function $\varepsilon: E \rightarrow(0, \infty)$ there exists a $C^{k}$ mapping $g: E \rightarrow F$ such that $\|f(x)-g(x)\| \leq \varepsilon(x)$ and $D g(x): E \rightarrow F$ is a surjective linear operator for every $x \in E$.

Here $k$ denotes the order of smoothness of the space $E$, defined as follows: $k=\infty$ if $E \in\left\{c_{0}\right\} \cup\left\{\ell_{2 n}\right.$ : $n \in \mathbb{N}\} \cup\left\{L^{2 n}: n \in \mathbb{N}\right\} ; k=2 n+1$ if $E \in\left\{\ell_{2 n+1}: n \in \mathbb{N}\right\} \cup\left\{L^{2 n+1}: n \in \mathbb{N}\right\}$, and $k$ is equal to the integer part of $p$ if $E \in\left\{\ell_{p}\right\} \cup\left\{L^{p}\right\}$ and $p \notin \mathbb{N}$. The Sobolev spaces $W^{k, p}\left(\mathbb{R}^{n}\right)$ with $1<p<\infty$ are also included in Theorem 1.2 since they are isomorphic to $L^{p}\left(\mathbb{R}^{n}\right)$ (see [56, Theorem 11]).
Notice that the assumption that there exists a bounded linear operator from $E$ onto $F$ is necessary, as otherwise all points of $E$ are critical for all functions $g \in C^{1}(E, F)$.
Of course Theorem 1.1 is a particular case of Theorem 1.2 (also note that if $E$ is a separable Hilbert space, $F$ is a Banach space, and there exists a continuous linear surjection $T: E \rightarrow F$, then $F$ must be isomorphic to $\mathbb{R}^{n}$ or to $E$ ). In general, note that a continuous linear surjection $T: E \rightarrow F$ between Banach spaces exists if and only if $F$ is isomorphic to a quotient space of $E$.

[^0]We will also establish more technical results (see Theorems 1.6 and 1.7 below) that generalize the preceding theorems to much larger classes of Banach spaces (especially in the case that $E$ is reflexive). Part of the motivation for this kind of results is in their connection with the Morse-Sard theorem, a fundamental result in Differential Geometry and Analysis. Throughout this paper, for a $C^{k}$ smooth mapping $f: \mathbb{R}^{n} \longrightarrow \mathbb{R}^{m}, C_{f}$ stands for the set of critical points of $f$ (that is, the points $x \in \mathbb{R}^{n}$ at which the differential $D f(x)$ is not surjective), and $f\left(C_{f}\right)$ is thus the set of critical values of $f$; the same terminology applies to smooth mappings between manifolds, both finite and infinite-dimensional. The Morse-Sard theorem [54, 62] states that if $k \geq \max \{n-m+1,1\}$ then $f\left(C_{f}\right)$ is of Lebesgue measure zero in $\mathbb{R}^{m}$. This result also holds true for $C^{k}$ smooth mappings $f: N \longrightarrow M$ between two smooth manifolds of dimensions $n$ and $m$ respectively.
Given the crucial applications of the Morse-Sard theorem in several branches of mathematics, it is natural both to try to extend this result for other classes of mappings, and also to ask what happens in the case that $M$ and $N$ are infinite-dimensional manifolds. Regarding the first issue, many refinements of the Morse-Sard theorem for other classes of mappings (notably Hölder, Sobolev, and BV mappings) have appeared in the literature; see for instance [68, 69, $55,11,12, ~ 53, ~ 24, ~ 35, ~ 18, ~ 19, ~ 45, ~ 41, ~ 40, ~ 9, ~ 10] ~$ and the references therein.
As for the second issue, which in this paper is of our concern, let us mention the results of several authors who have studied the question as to what extent one can obtain results similar to the MorseSard theorem for mappings between infinite-dimensional Banach spaces or manifolds modeled on such spaces.
S. Smale [63] proved that if $X$ and $Y$ are separable connected smooth manifolds modeled on Banach spaces and $f: X \longrightarrow Y$ is a $C^{r}$ Fredholm mapping (that is, every differential $D f(x)$ is a Fredholm operator between the corresponding tangent spaces) then $f\left(C_{f}\right)$ is meager, and in particular $f\left(C_{f}\right)$ has no interior points, provided that $r>\max \{\operatorname{index}(D f(x)), 0\}$ for all $x \in X$; here index $(D f(x))$ stands for the index of the Fredholm operator $\operatorname{Df}(x)$, that is, the difference between the dimension of the kernel of $D f(x)$ and the codimension of the image of $D f(x)$, both of which are finite. Of course, these assumptions are very restrictive as, for instance, if $X$ is infinite-dimensional then no function $f: X \longrightarrow \mathbb{R}$ is Fredholm.
In general, every attempt to adapt the Morse-Sard theorem to infinite dimensions will have to impose vast restrictions because, as shown by Kupka's counterexample [47], there are $C^{\infty}$ smooth functions $f: \ell_{2} \longrightarrow \mathbb{R}$ so that their sets of critical values $f\left(C_{f}\right)$ contain intervals. Furthermore, as shown by Bates and Moreira in [12, 53, one can take $f$ to be a polynomial of degree 3.
Nevertheless, for many applications of the Morse-Sard theorem, it is often enough to know that any given continuous mapping can be uniformly approximated by a mapping whose set of critical values is small in some sense; therefore it is natural to ask what mappings between infinite-dimensional manifolds will at least have such an approximation property. Going in this direction, Eells and McAlpin established the following theorem [32]: If $E$ is a separable Hilbert space, then every continuous function from $E$ into $\mathbb{R}$ can be uniformly approximated by a smooth function $f$ whose set of critical values $f\left(C_{f}\right)$ is of measure zero. This allowed them to deduce a version of this theorem for mappings between smooth manifolds $M$ and $N$ modeled on $E$ and a Banach space $F$ respectively, which they called an approximate Morse-Sard theorem: Every continuous mapping from $M$ into $N$ can be uniformly approximated by a smooth mapping $f: M \longrightarrow N$ so that $f\left(C_{f}\right)$ has empty interior. However, as observed in [32, Remark 3A], we have $C_{f}=M$ in the case that $F$ is infinite-dimensional (so, even though the set of critical values of $f$ is relatively small, the set of critical points of $f$ is huge, which is somewhat disappointing).
In [6], a much stronger result was obtained by M. Cepedello-Boiso and the first-named author: if $M$ is a $C^{\infty}$ smooth manifold modeled on a separable infinite-dimensional Hilbert space $X$, then every continuous mapping from $M$ into $\mathbb{R}^{m}$ can be uniformly approximated by smooth mappings with no
critical points. P. Hájek and M. Johanis [37] established a similar result for $m=1$ in the case that $X$ is a separable Banach space which contains $c_{0}$ and admits a $C^{p}$-smooth bump function. Finally, in the case that $m=1$, these results were extended by M. Jiménez-Sevilla and the first-named author [7] for continuous functions $f: X \rightarrow \mathbb{R}$, where $X$ is a separable Banach space admitting an equivalent smooth and locally uniformly rotund norm.
In this paper, we will improve these results by showing that the pairs $\left(\ell_{2}, \mathbb{R}^{m}\right)$ or $(X, \mathbb{R})$ can be replaced with pairs of the form $(E, F)$, where $E$ is a Banach space from a large class (including all the classical spaces with smooth norms such as $c_{0}, \ell_{p}$ or $\left.L^{p}, 1<p<\infty\right)$, and $F$ can be taken to be any quotient space of $E$. So we may say that even though an exact Morse-Sard theorem for mappings between classical Banach spaces is false, a stronger approximate version of the Morse-Sard theorem is nonetheless true.
The general plan of the proof of Theorem 1.2 consists in following these steps:

- Step 1: We construct a smooth mapping $\varphi: E \rightarrow F$ such that $\|\varphi(x)-f(x)\| \leq \varepsilon(x) / 2$ and $C_{\varphi}$, the critical set of $\varphi$, is locally contained in the graph of a continuous mapping defined on a complemented subspace of infinite codimension in $E$ and taking values in its linear complement.
- Step 2: We find a diffeomorphism $h: E \rightarrow E \backslash C_{\varphi}$ such that $h$ is sufficiently close to the identity, in the sense that $\{\{x, h(x)\}: x \in E\}$ refines $\mathcal{G}$ (in other words, $h$ is limited by $\mathcal{G}$ ), where $\mathcal{G}$ is an open cover of $E$ by open balls $B\left(z, \delta_{z}\right)$ chosen in such a way that if $x, y \in B\left(z, \delta_{z}\right)$ then

$$
\|\varphi(y)-\varphi(x)\| \leq \frac{\varepsilon(z)}{4} \leq \frac{\varepsilon(x)}{2} .
$$

The existence of such a diffeomorphism $h$ follows by the results of Section 2.

- Step 3: Then, the mapping $g(x):=\varphi(h(x))$ has no critical point and satisfies $\|f(x)-g(x)\| \leq$ $\varepsilon(x)$ for all $x \in E$.
The results of Section 2 are of independent interest, as they generalize important theorems on diffeomorphic extractions of some kind of sets. Although it is well known (see [20, 51, 30, 31] and the references therein) that every two separable, homotopy equivalent, infinite-dimensional Hilbert manifolds $M, N$ are in fact diffeomorphic, a diffeomorphism $h: M \rightarrow N$ provided by this deep result has not been (and, in general, cannot be) shown to be limited by an arbitrary open cover $\mathcal{G}$ of $M$, a property that is essential in Step 3 above. The finest result we know of which provides a diffeomorphism $h: E \rightarrow E \backslash X$ limited by a given open cover $\mathcal{G}$ of $E$, where $E$ is a separable infinite-dimensional Hilbert space $E$ and $X$ is a closed subset of $E$, is a theorem of J.E. West [67] in which $X$ is assumed to be locally compact. However, in the proof of Theorem [1.2, we do not work necessarily with Hilbert spaces and we need to diffeomorphically extract a closed set $X$ which is not necessarily locally compact but merely locally contained in the graph of a continuous mapping defined on a complemented subspace of infinite codimension in $E$ and taking values in its linear complement (for a precise explanation of this terminology, see the statement of Theorem 1.4 below). In Section 2, we construct diffeomorphisms $h$ which extract such closed sets $X$.
The main result of Section 2 is the following.
Theorem 1.3. Let $E$ be an infinite-dimensional Hilbert space, $X$ a closed subset of $E$ which is locally contained in the graph of a continuous function defined on a subspace of infinite codimension in $E$ and taking values in its orthogonal complement, $\mathcal{G}$ an open cover of $E$, and $U$ an open subset of $E$. Then, there exists a $C^{\infty}$ diffeomorphism $h$ of $E \backslash X$ onto $E \backslash(X \backslash U)$ which is the identity on $(E \backslash U) \backslash X$ and is limited by $\mathcal{G}$.
Recall that $h$ is said to be limited by $\mathcal{G}$ provided that the set $\{\{x, h(x)\}: x \in E \backslash X\}$ refines $\mathcal{G}$; that is, for every $x \in E \backslash X$, we may find a $G_{x} \in \mathcal{G}$ such that both $x$ and $h(x)$ are in $G_{x}$.
Theorem 1.3 is a straightforward consequence of the following much more general result, which is true for many Banach spaces not necessarily Hilbertian.

Theorem 1.4. Let $E$ be a Banach space, $p \in \mathbb{N} \cup\{\infty\}$, and $X \subset E$ be a closed set with the property that, for each $x \in X$, there exist a neighborhood $U_{x}$ of $x$ in $E$, Banach spaces $E_{(1, x)}$ and $E_{(2, x)}$, and a continuous mapping $f_{x}: C_{x} \rightarrow E_{(2, x)}$, where $C_{x}$ is a closed subset of $E_{(1, x)}$, such that:
(1) $E=E_{(1, x)} \oplus E_{(2, x)}$;
(2) $E_{(1, x)}$ has $C^{p}$ smooth partitions of unity;
(3) $E_{(2, x)}$ is infinite-dimensional and has a (not necessarily equivalent) norm of class $C^{p}$;
(4) $X \cap U_{x} \subset G\left(f_{x}\right)$, where

$$
G\left(f_{x}\right)=\left\{y=\left(y_{1}, y_{2}\right) \in E_{(1, x)} \oplus E_{(2, x)}: y_{2}=f_{x}\left(y_{1}\right), y_{1} \in C_{x}\right\} .
$$

Then, for every open cover $\mathcal{G}$ of $E$ and every open subset $U$ of $E$, there exists a $C^{p}$ diffeomorphism $h$ from $E \backslash X$ onto $E \backslash(X \backslash U)$ which is the identity on $(E \backslash U) \backslash X$ and is limited by $\mathcal{G}$. Moreover, the same conclusion is true if we replace $E$ with an open subset of $E$.

The proof of Theorem 1.4 combines ideas and techniques from Peter Renz's Ph.D. thesis [58], James West's paper [67, and some previous work of the first and second-named authors [4, 5]; see Section 2 for more information. It should be noted that (even in the case that $E$ is separable) Theorem 1.3 generalizes West's theorem [67], because a closed locally compact subset of an infinite-dimensional Hilbert space $E$, locally, can be regarded as the graph of a continuous mapping defined on a closed subset of an infinite-codimensional subspace of $E$; see, for instance, 58]. Furthermore, note in the above results we do not assume separability of the Banach space $E$.
The proof of Theorem 1.1 will show that $E$ and $F$ can be replaced with open subsets $U$ and $V$ of $E$ and $F$ respectively. Then, by combining such an equivalent statement of Theorem 1.1 with the well known result [30, 46] stating that every separable infinite-dimensional Hilbert manifold is diffeomorphic to an open subset of $\ell_{2}$, one may easily deduce the following.

Theorem 1.5. Let $M, N$ be separable infinite-dimensional Hilbert manifolds. For every continuous mapping $f: M \rightarrow N$ and every open cover $\mathcal{U}$ of $N$, there exists a $C^{\infty}$ mapping $g: M \rightarrow N$ such that $g$ has no critical point and $\{\{f(x), g(x)\}: x \in M\}$ refines $\mathcal{U}$.

Alternatively, one can also adjust the proof of Theorem 1.1 to obtain a direct proof of Theorem 1.5
It is worth noting that Theorems 1.1 and 1.2 are immediate consequences of the following more general (but also more technical) results. For spaces $E$ which are reflexive and have a certain "composite" structure, we have the following.

Theorem 1.6. Let $E$ be a separable reflexive Banach space of infinite dimension, and $F$ be a Banach space. In the case that $F$ is infinite-dimensional, let us assume furthermore that:
(1) $E$ is isomorphic to $E \oplus E$.
(2) There exists a linear bounded operator from $E$ onto $F$ (equivalently, $F$ is a quotient space of E).

Then, for every continuous mapping $f: E \rightarrow F$ and every continuous function $\varepsilon: E \rightarrow(0, \infty)$ there exists a $C^{1}$ mapping $g: E \rightarrow F$ such that $\|f(x)-g(x)\| \leq \varepsilon(x)$ and $D g(x): E \rightarrow F$ is a surjective linear operator for every $x \in E$.

Note that there exists separable, reflexive Banach spaces $E$ such that $E$ is not isomorphic to $E \oplus E$. The first example of such a space was given by Figiel in 1972 [34].
See also Theorem 5.1 and Theorem 5.2 below for more general variants of this result.
For spaces which are not necessarily reflexive but have an appropriate Schauder basis we have the following.

Theorem 1.7. Let $E$ be an infinite-dimensional Banach space, and $F$ be a Banach space such that:
(1) $E$ has an equivalent locally uniformly convex norm $\|\cdot\|$ which is $C^{1}$ smooth.
(2) $E=(E,\|\cdot\|)$ has a (normalized) Schauder basis $\left\{e_{n}\right\}_{n \in \mathbb{N}}$ such that for every $x=\sum_{j=1}^{\infty} x_{j} e_{j}$ and every $j_{0} \in \mathbb{N}$ we have that

$$
\left\|\sum_{j \in \mathbb{N}, j \neq j_{0}} x_{j} e_{j}\right\| \leq\left\|\sum_{j \in \mathbb{N}} x_{j} e_{j}\right\|
$$

(3) In the case that $F$ is infinite-dimensional, there exists a subset $\mathbb{P}$ of $\mathbb{N}$ such that both $\mathbb{P}$ and $\mathbb{N} \backslash \mathbb{P}$ are infinite and, for every infinite subset $J$ of $\mathbb{P}$, there exists a linear bounded operator

Then, for every continuous mapping $f: E \rightarrow F$ and every continuous function $\varepsilon: E \rightarrow(0, \infty)$ there exists a $C^{1}$ mapping $g: E \rightarrow F$ such that $\|f(x)-g(x)\| \leq \varepsilon(x)$ and $D g(x): E \rightarrow F$ is a surjective linear operator for every $x \in E$.

Recall that a norm $\|\cdot\|$ in a Banach space $E$ is said to be locally uniformly convex (LUC) (or locally uniformly rotund (LUR)) provided that, for every sequence $\left(x_{n}\right) \subset E$ and every point $x_{0}$ in $E$, we have that

$$
\lim _{n \rightarrow \infty} 2\left(\left\|x_{0}\right\|^{2}+\left\|x_{n}\right\|^{2}\right)-\left\|x_{0}+x_{n}\right\|^{2}=0 \Longrightarrow \lim _{n \rightarrow \infty}\left\|x_{n}-x_{0}\right\|=0
$$

Condition (2) is equivalent to the fact that, for every (equivalently, finite) set $A \subset \mathbb{N},\left\|P_{A}\right\| \leq 1$, where $P_{A}$ stands for the projection $P_{A}(x)=\sum_{j \in A} x_{j} e_{j}$. This, in particular, implies that $\left\{e_{n}\right\}_{n \in N}$ is an unconditional basis; for more details see [2, p. 53] or [1].
The proofs of these theorems will be provided in Sections 3 and 4. These results combine to yield Theorem 1.2 for $k=1$ (see also Remark 5.5 in Section 5 for an explanation of why the space $c_{0}$ satisfies the assumptions of Theorem [1.7). In order to deduce Theorem [1.2 in the cases of higher order smoothness, we just have to use Nicole Moulis's results on $C^{1}$ fine approximation in Banach spaces [52] or the more general results of [38, Corollary 7.96], together with the following fact.

Proposition 1.8. Assume that the Banach spaces E, F satisfy the following properties:
(1) For every continuous mapping $f: E \rightarrow F$ and every continuous function $\delta: E \rightarrow(0, \infty)$ there exists a $C^{1}$ smooth mapping $\varphi: E \rightarrow F$ such that $\|f(x)-\varphi(x)\| \leq \delta(x)$ and $D \varphi(x): E \rightarrow F$ is surjective for all $x \in E$.
(2) For every $C^{1}$ mapping $\varphi: E \rightarrow F$ and every continuous function $\eta: E \rightarrow(0, \infty)$ there exists a $C^{k}$ mapping $g: E \rightarrow F$ such that $\|f(x)-\varphi(x)\| \leq \eta(x)$ and $\|D \varphi(x)-D g(x)\| \leq \eta(x)$ for all $x \in E$.
Then, for every continuous mapping $f: E \rightarrow F$ and every continuous function $\varepsilon: E \rightarrow(0, \infty)$ there exists a $C^{k}$ smooth mapping $g: E \rightarrow F$ such that $\|f(x)-g(x)\| \leq \varepsilon(x)$ and $D g(x): E \rightarrow F$ is surjective for every $x \in E$.

Nevertheless, it should be noted that our proof of Theorem 1.6 directly provides $C^{\infty}$ approximations without critical points in the case that $E$ is a separable Hilbert space; see Remark 5.3 in Section 5 below. An easy proof of Proposition 1.8, together with some examples, remarks and more technical variants of our results, is given in Section 5 .
Finally, let us mention that as a straightforward application of Theorem [1.2, we obtain that, for all Banach spaces $E$ and $F$ appearing in Theorem [1.2, every continuous mapping $f: E \rightarrow F$ can be uniformly approximated by open mappings of class $C^{k}$. For a more general statement, see Remark 5.9. Obviously, the latter result is false in the case that $E$ is finite-dimensional.

## 2. EXtracting closed sets which are locally contained in graphs of infinite CODIMENSION

In this section we will combine ideas and tools of [58, 67, 5] in order to prove Theorem 1.4 , We will split the proof into four subsections. First, in Section 2.1, we will see that each piece of $X$ contained in the graph $G\left(f_{x}\right)$ as provided by condition (4) of the statement can be flattened by means of homeomorphisms $h_{x}, \varphi_{x}: E \rightarrow E$ which are sufficiently close to each other, and whose restrictions to $E \backslash G\left(f_{x}\right)$ and $E \backslash\left(G\left(f_{x}\right) \backslash U\right)$ are diffeomorphisms, respectively. Next, in Sections 2.2 and 2.3, we will show that there exists a diffeomorphism $g_{x}: E \backslash\left(C_{x} \times\{0\}\right) \rightarrow E \backslash\left(\left(C_{x} \times\{0\}\right) \backslash h_{x}(U)\right)$ which is the identity on $\left(E \backslash h_{x}(U)\right) \backslash\left(C_{x} \times\{0\}\right)$ and moves no point more than a fixed small number $\varepsilon_{x}$. Then, the composition $\varphi_{x}^{-1} \circ g_{x} \circ h_{x}$ will extract the local chunk of graph $U_{x} \cap G\left(f_{x}\right)$ and will move no point too much. Finally, in Section 2.4 , we will see how one can patch a collection of diffeomorphisms extracting pieces of $X$ into a diffeomorphism $h$ which extracts $X$ and is limited by $\mathcal{G}$.
In Section 2.1, we will closely follow Peter Renz's results from [58, 59]. In Sections 2.2 and 2.3, we will combine ideas and techniques from [58, 4, 5]. Finally, in Section 2.4, we will borrow a technique of James West's [67, p. 288-290].
2.1. Flattening graphs. Here we will prove the following.

Theorem 2.1. Let $E_{1}$ be a Banach space with $C^{p}$ smooth partitions of unity and $E_{2}$ be a Banach space which admits a (not necessarily equivalent) $C^{p}$ norm. Let $\left(E=E_{1} \times E_{2},\|\cdot\|\right)$ and $\pi_{1}: E \rightarrow E_{1}$ be the natural projection, i.e., $\pi_{1}\left(x_{1}, x_{2}\right)=x_{1},\left(x_{1}, x_{2}\right) \in E$. Let $X_{1} \subset E_{1}$ be a closed set, $f: X_{1} \rightarrow E_{2} a$ continuous mapping, $U \subset E$ an open set, and $\varepsilon>0$. Write $G(f)=\left\{\left(x_{1}, x_{2}\right) \in E: x_{2}=f\left(x_{1}\right)\right.$, $x_{1} \in$ $\left.X_{1}\right\}$. Then there exist a couple of homeomorphisms $h, \varphi: E \rightarrow E$ such that:
(1) $h(G(f)) \subset E_{1} \times\{0\}$ and $\varphi(G(f) \backslash U) \subset\left(E_{1} \times\{0\}\right) \backslash h(U)$;
(2) $h=\varphi$ off of $U$;
(3) $\pi_{1} \circ h=\pi_{1}=\pi_{1} \circ \varphi$;
(4) $h$ restricted to $E \backslash G(f)$ is a $C^{p}$ diffeomorphism of $E \backslash G(f)$ onto $E \backslash\left(X_{1} \times\{0\}\right)$;
(5) $\varphi$ restricted to $E \backslash(G(f) \backslash U)$ is a $C^{p}$ diffeomorphism of $E \backslash(G(f) \backslash U)$ onto $E \backslash\left(\left(X_{1} \times\{0\}\right) \backslash h(U)\right)$.
(6) $\left\|h^{-1}(x)-\varphi^{-1}(x)\right\| \leq \varepsilon$ for every $x \in E$.
(7) $h^{-1}\left(x_{1}, x_{2}\right)$ is uniformly continuous with respect to the second coordinate $x_{2}$ (meaning that for every $\varepsilon>0$ there exists $\delta>0$ such that if $\left\|x_{2}-x_{2}^{\prime}\right\|<\delta$ then $\left\|h^{-1}\left(x_{1}, x_{2}\right)-h^{-1}\left(x_{1}, x_{2}^{\prime}\right)\right\|<\varepsilon$ for all $x_{1}$ ).

We will assume without loss of generality that $\varepsilon \leq 1$.
In what follows, slightly abusing notation, we will indistinctly use the symbol $\|\cdot\|$ to denote the norms $\|\cdot\|_{E_{1}},\|\cdot\|_{E_{2}}$, and $\|\cdot\|$ with which the Banach spaces $E_{1}, E_{2}$ or $E_{1} \times E_{2}$ are endowed. We may and do assume that $\left\|x_{1}\right\|_{E_{1}}=\left\|\left(x_{1}, 0\right)\right\|$ and $\left\|x_{2}\right\|_{E_{2}}=\left\|\left(0, x_{2}\right)\right\|$ for all $\left(x_{1}, x_{2}\right) \in E_{1} \times E_{2}$.
Now, we state and prove a sequence of lemmas that will be employed in proving the above theorem. The most important are Lemmas 2.3 and [2.5. Basically, we follow the ideas of Renz's paper [59] and Ph.D. thesis [58], with some minor but very important changes.
The proof of our first lemma is a consequence of the existence of $C^{p}$ smooth partitions of unity on $E_{1}$.
Lemma 2.2. The function $f: X_{1} \rightarrow E_{2}$ extends to a continuous function $\bar{f}: E_{1} \rightarrow E_{2}$ such that $\bar{f} \mid E_{1} \backslash X_{1}$ is $C^{p}$ smooth.
Lemma 2.3. Let $E_{1}$ be a Banach space with $C^{p}$ smooth partitions of unity, $E_{2}$ be a Banach space removed text, $X_{1}$ be a closed subset of $E_{1}$, and $f: X_{1} \rightarrow E_{2}$ be a continuous mapping. For every $n \in \mathbb{N}$, write

$$
W_{n}=\left\{x_{1} \in E_{1}: \operatorname{dist}\left(x_{1}, X_{1}\right) \leq \frac{1}{n}\right\}
$$

Assume $\bar{f}: E_{1} \rightarrow E_{2}$ is a continuous extension of $f$ such that $\bar{f} \mid E_{1} \backslash X_{1}$ is $C^{p}$ smooth. Then, there is a continuous mapping

$$
F: \mathbb{R} \times E_{1} \rightarrow E_{2}
$$

such that
(1) $F\left(r, x_{1}\right)=\bar{f}\left(x_{1}\right)$ for all $\left(r, x_{1}\right) \in\left(r_{n}, \infty\right) \times E_{1} \backslash W_{n}$ and some $0<r_{n}<1$; in particular, $F\left(r, x_{1}\right)=f\left(x_{1}\right)$ for all $\left(r, x_{1}\right)$ in some neighborhood of the set $\{1\} \times\left(E_{1} \backslash X_{1}\right)$ in $\mathbb{R} \times\left(E_{1} \backslash X_{1}\right)$;
(2) $F \mid \mathbb{R} \times\left(E_{1} \backslash X_{1}\right) \cup(-\infty, 1) \times E_{1}$ is $C^{p}$ smooth;
(3) $F\left(r, x_{1}\right)=f\left(x_{1}\right)$ for $r \geq 1$ and $x_{1} \in X_{1}$;
(4) $\left\|D_{1} F\left(r, x_{1}\right)\right\| \leq \frac{1}{2}$ for all $r \in \mathbb{R}, x_{1} \in E_{1}$.

Proof. For every $n \in \mathbb{N}$ we can find a sequence of $C^{p}$ functions $\bar{f}_{n}: E_{1} \rightarrow E_{2}$ such that

$$
\left\|\bar{f}\left(x_{1}\right)-\bar{f}_{n}\left(x_{1}\right)\right\| \leq 2^{-2 n-4}
$$

for every $x_{1} \in E_{1}$. The existence of such a sequence is again guaranteed by the existence of $C^{p}$ partitions of unity in $E_{1}$ (see, for instance, [26, Theorem VII.3.2]). We will now improve the sequence $\left\{\bar{f}_{n}\right\}_{n \geq 1}$ to $\left\{f_{n}\right\}_{n \geq 1}$ so that the sequence $\left\{f_{n} \mid E_{1} \backslash X_{1}\right\}$ locally stabilizes with respect to $n$. To achieve this, we use the existence of $C^{p}$ partitions of unity to find a $C^{p}$ function $\lambda_{n}: E_{1} \rightarrow[0,1]$ which is 1 on $E \backslash W_{n}$ and 0 on $W_{n+1}$. Define

$$
f_{n}\left(x_{1}\right)=\lambda_{n}\left(x_{1}\right) \bar{f}\left(x_{1}\right)+\left(1-\lambda_{n}\left(x_{1}\right)\right) \bar{f}_{n}\left(x_{1}\right)
$$

for all $x_{1} \in E_{1}$. It follows that $\left\|f_{n}\left(x_{1}\right)-f_{n+1}\left(x_{1}\right)\right\| \leq 2^{-2 n-3}$ for all $x_{1} \in E_{1}$.
For every $n \in \mathbb{N}$, pick a nondecreasing $C^{\infty}$ function $h_{n}: \mathbb{R} \rightarrow[0,1]$ such that $h_{n}(r)=0$ for $r \leq 1-2^{1-n}$, $h_{n}(r)=1$ for $r \geq 1-2^{-n}$, and $h_{n}^{\prime}(r) \leq 2^{n+1}$. One can check that

$$
F\left(r, x_{1}\right)=f_{1}\left(x_{1}\right)+\sum_{n=1}^{\infty} h_{n+1}(r)\left(f_{n+1}\left(x_{1}\right)-f_{n}\left(x_{1}\right)\right)
$$

defines a required mapping.
Observe that in fact $F\left(r, x_{1}\right)$ is Lipschitz with constant 1 with respect to the first variable $r \in \mathbb{R}$. That is,

$$
\begin{aligned}
\left\|F\left(r, x_{1}\right)-F\left(r^{\prime}, x_{1}\right)\right\| & \leq \sum_{n=1}^{\infty}\left|h_{n+1}(r)-h_{n+1}\left(r^{\prime}\right)\right|| | f_{n+1}\left(x_{1}\right)-f_{n}\left(x_{1}\right) \| \leq \\
& \leq \sum_{n=1}^{\infty} 2^{n+1}\left|r-r^{\prime}\right| 2^{-2 n-3} \leq\left|r-r^{\prime}\right|
\end{aligned}
$$

for every $x_{1} \in E_{1}$.
We will write

$$
U_{0}=\pi_{1}(G(\bar{f}) \cap U),
$$

which is an open set in $E_{1}$, and also

$$
Y_{1}=X_{1} \backslash U_{0}=X_{1} \backslash \pi_{1}(G(\bar{f}) \cap U)=\pi_{1}(G(f) \backslash U)
$$

which is a closed subset of $E_{1}$. By replacing $U$ with $U \cap \pi_{1}^{-1}\left(U_{0}\right)$, we can assume that

$$
U_{0}=\pi_{1}(U)
$$

Lemma 2.4. With the above notation, take a decreasing sequence of positive numbers $\left\{\delta_{n}\right\}_{n \geq 1}$ converging to zero. Then there exists an increasing sequence of open subsets in $E_{1}$

$$
V_{1} \subseteq \overline{V_{1}} \subseteq V_{2} \subseteq \overline{V_{2}} \subseteq \cdots \subseteq \overline{V_{n}} \subseteq V_{n+1} \subseteq \cdots \subseteq U_{0}
$$

such that $\bigcup_{n=1}^{\infty} V_{n}=U_{0}$ and the sets

$$
U_{n}:=\left\{\left(x_{1}, x_{2}\right) \in E:\left\|x_{2}-\bar{f}\left(x_{1}\right)\right\|<\delta_{n}, x_{1} \in V_{n}\right\}
$$

are contained in $U$.
Proof. To be able to get the required inclusions between the sets $V_{n}$, we first take an auxiliary sequence of open sets $W_{n}$ in $U_{0}$ such that $\overline{W_{n}} \subseteq W_{n+1}$ for every $n \in N$ and $\bigcup_{n=1}^{\infty} W_{n}=U_{0}$.
Then we define

$$
V_{n}^{\prime}=\left\{x_{1} \in U_{0}:\left\{\left(x_{1}, x_{2}\right) \in E:\left\|x_{2}-\bar{f}\left(x_{1}\right)\right\|<\delta_{n}\right\} \subseteq U\right\}
$$

for every $n \in \mathbb{N}$. Observe that we have $V_{n}^{\prime} \subseteq V_{n+1}^{\prime}$ and $\bigcup_{n=1}^{\infty} V_{n}^{\prime}=U_{0}$, but we cannot assure that $\overline{V_{n}^{\prime}} \subseteq$ $V_{n+1}^{\prime}$ for every $n \in N$. So now we mix these sets with the previous $W_{n}$, that is, we let $V_{n}=W_{n} \cap V_{n}^{\prime}$. Obviously, by definition, for every $n \in \mathbb{N}$ the set $U_{n}=\left\{\left(x_{1}, x_{2}\right) \in E:\left\|x_{2}-\bar{f}\left(x_{1}\right)\right\|<\delta_{n}, x_{1} \in V_{n}\right\}$ is contained in $U$. Now, we have that $\overline{V_{n}} \subseteq V_{n+1}$ for every $n \in N$; also $\bigcup_{n=1}^{\infty} V_{n}=U_{0}$.

The following lemma resembles [59, Lemma 2.2] and [58, Lemma 2] (in which only one function $\phi$ is considered). However, Theorem [2.1] requires constructing two homeomorphisms $h$ and $\varphi$ which are identical outside $U$. The building block in constructing those homeomorphisms are two functions $\phi$ and $\tilde{\phi}$ whose existence is claimed in the lemma below. The existence of $\tilde{\phi}$ is crucial. Incidentally, let us note that Renz's proof of [58, Theorem 4] is flawed (and this is the reason why we must deal with two functions $\phi$ and $\tilde{\phi}$ instead of just the function $\phi$ ), but can be corrected by using Theorem 2.1.

Lemma 2.5. Let $\bar{f}: E_{1} \rightarrow E_{2}$ be the uniform limit of $C^{p}$ functions, where $E_{1}$ has $C^{p}$ partitions of unity and $E_{2}$ has a (not necessarily equivalent) $C^{p}$ smooth norm. Then there are two continuous functions $\phi, \tilde{\phi}: E \rightarrow[0,1]$ such that
(1) $\phi^{-1}(1)=G(\bar{f})$ and $\tilde{\phi}^{-1}(1)=G(\bar{f}) \backslash U$;
(2) $\phi \mid E \backslash G(\bar{f})$ and $\tilde{\phi} \mid E \backslash(G(\bar{f}) \backslash U)$ are $C^{p}$ smooth;
(3) $\left\|D_{2} \phi\left(x_{1}, x_{2}\right)\right\| \leq \frac{1}{2}$ for all $\left(x_{1}, x_{2}\right) \in E \backslash G(\bar{f})$, and $\left\|D_{2} \tilde{\phi}\left(x_{1}, x_{2}\right)\right\| \leq \frac{1}{2}$ for all $\left(x_{1}, x_{2}\right) \in$ $E \backslash(\underset{\sim}{G}(\bar{f}) \backslash U) ;$
(4) $\phi=\tilde{\phi}$ outside $U$.

Proof. To construct $\phi$ we will follow [59, Lemma 2.2]. A similar argument will be used to construct $\tilde{\phi}$; however, we have to make sure that $\tilde{\phi} \mid G(\bar{f}) \cap U<1$.
For $n \in \mathbb{N}$, let $a_{n}, b_{n}, c_{n}, d_{n}, \varepsilon_{n}$ be positive numbers with the following properties:
(1) they tend to zero as $n$ tends to infinity;
(2) $a_{n}<b_{n}$ for all $n$;
(3) $\varepsilon_{n+1}+b_{n+1}<a_{n}-\varepsilon_{n}$ for all $n$;
(4) $\sum_{n=1}^{\infty} c_{n} \leq \frac{\epsilon}{2} \leq 1$;
(5) $\sum_{n=1}^{\infty} d_{n} \leq \frac{1}{2}$
(for instance, let us set $a_{n}=\epsilon 2^{-2 n}, b_{n}=\epsilon 2 a_{n}, c_{n}=\epsilon 2^{-4 n}, d_{n}=\epsilon 2^{-2 n}$ and $\varepsilon_{n}=\epsilon 2^{-4(n+1)}$ ). Let $h_{n}$ be a nonincreasing $C^{p}$ function from $\mathbb{R}$ to $\mathbb{R}$ satisfying

$$
\begin{array}{lr}
c_{n}=h_{n}(r)=h_{n}(0)>0 & \text { whenever } r \leq a_{n}, \\
h_{n}(r)=0 & \text { whenever } r \geq b_{n}, \\
\left|h_{n}^{\prime}(r)\right| \leq d_{n} & \text { for all } r \text { in } \mathbb{R},
\end{array}
$$

and $g_{n}: E_{1} \rightarrow E_{2}$ be a $C^{p}$ mapping such that $\left\|g_{n}\left(x_{1}\right)-\bar{f}\left(x_{1}\right)\right\| \leq \varepsilon_{n}$ for every $x_{1} \in E_{1}$. Then

$$
\psi_{n}\left(x_{1}, x_{2}\right)=h_{n}\left(\left\|x_{2}-g_{n}\left(x_{1}\right)\right\|\right)
$$

defines a nonnegative $C^{p}$ function on $E_{1} \times E_{2}=E$ satisfying

$$
\begin{array}{ll}
c_{n}=\psi_{n}\left(x_{1}, x_{2}\right)=h_{n}(0)>0 & \text { if }\left\|x_{2}-\bar{f}\left(x_{1}\right)\right\| \leq a_{n}-\varepsilon_{n} \\
\psi_{n}\left(x_{1}, x_{2}\right)=0 & \text { if }\left\|x_{2}-\bar{f}\left(x_{1}\right)\right\| \geq b_{n}+\varepsilon_{n} \\
\left\|D_{2} \psi_{n}\left(x_{1}, x_{2}\right)\right\| \leq d_{n} & \text { for all }\left(x_{1}, x_{2}\right) \in E_{1} \times E_{2}
\end{array}
$$

The nonnegativity and first two properties of $\psi_{n}$ are evident, and it is easy to see that $\psi_{n}$ is $C^{p}$ on $E$. The bound on the norm of the derivative $D_{2} \psi_{n}$ is established by using the chain rule and the fact that the operator norm of the derivative of the norm of any Banach space is less than or equal to one. Define

$$
\begin{equation*}
\psi\left(x_{1}, x_{2}\right)=\sum_{n=1}^{\infty} \psi_{n}\left(x_{1}, x_{2}\right) \tag{2.1}
\end{equation*}
$$

for all $\left(x_{1}, x_{2}\right) \in E$.
Similarly, define

$$
\begin{equation*}
\tilde{\psi}\left(x_{1}, x_{2}\right)=\sum_{n=1}^{\infty} \lambda_{n}\left(x_{1}\right) \psi_{n}\left(x_{1}, x_{2}\right) \tag{2.2}
\end{equation*}
$$

for all $\left(x_{1}, x_{2}\right) \in E$, where $\lambda_{n}: E_{1} \rightarrow[0,1]$ is a $C^{p}$ smooth function such that $\lambda_{n}\left(x_{1}\right)=1$ if $x_{1} \notin V_{n}$ and $\lambda_{n}\left(x_{1}\right)=0$ if $x \in V_{n-1}$. Here, the sets $V_{n}$ are provided by Lemma 2.4 for the sequence $\delta_{n}:=\varepsilon_{n}+b_{n}$ (let $V_{0}=\emptyset$ and assume $V_{1} \neq \emptyset$ ). In particular, observe that since $\bigcup_{n=1}^{\infty} V_{n}=U_{0}$ then $\lambda_{n}\left(x_{1}\right)=1$ for every $x_{1} \notin U_{0}$; hence, $\psi\left(x_{1}, x_{2}\right)=\tilde{\psi}\left(x_{1}, x_{2}\right)$ for $x_{1} \notin U_{0}$.
Since the functions $\psi$ and $\tilde{\psi}$ are defined via absolutely and uniformly convergent series of continuous functions, they are continuous.
If $\left(x_{1}, x_{2}\right) \in E \backslash G(\bar{f})$, then $\left\|x_{2}-\bar{f}\left(x_{1}\right)\right\|>b_{n}+\varepsilon_{n}$ for some $n \in \mathbb{N}$. By continuity, the inequality holds in a neighborhood of $\left(x_{1}, x_{2}\right)$ so $\psi_{k}$ vanishes for $k \geq n$. Hence, $\psi$ is locally a finite sum of $C^{p}$ functions, and in particular is $C^{p}$ on $E \backslash G(\bar{f})$.
Also, if $\left(x_{1}, x_{2}\right) \in \underset{\sim}{G}(\bar{f}) \cap U$, then $x_{1} \in U_{0}$ and $x_{1} \in V_{n}$ for some $n \in \mathbb{N}$. So $\lambda_{k}\left(x_{1}\right)=0$ for all $k \geq n+1$. This means that $\tilde{\psi}$ is locally a finite sum of $C^{p}$ functions and, thus, is of class $C^{p}$ on $E \backslash(G(\bar{f}) \backslash U)$. The derived series for $D_{2} \psi$ and $D_{2} \tilde{\psi}$ are absolutely and uniformly convergent in view of the bounds on $\left\|D_{2} \psi_{n}\right\|$ and the fact that $D_{2} \lambda_{n}=0$. Then differentiation term by term is justified and $\left\|D_{2} \psi\left(x_{1}, x_{2}\right)\right\| \leq$ $\frac{1}{2}$ for all $\left(x_{1}, x_{2}\right) \in E \backslash G(\bar{f})$ and $\left\|D_{2} \tilde{\psi}\left(x_{1}, x_{2}\right)\right\| \leq \frac{1}{2}$ for all $\left(x_{1}, x_{2}\right) \in E \backslash(G(\bar{f}) \backslash U)$.
Each point $\left(x_{1}, x_{2}\right) \in G(\bar{f})$ satisfies $0=\left\|x_{2}-\bar{f}\left(x_{1}\right)\right\|<a_{n}-\varepsilon_{n}$ for all $n \in \mathbb{N}$, consequently $\psi$ equals the constant

$$
d^{*}=\sum_{n=1}^{\infty} \psi_{n}\left(x_{1}, \bar{f}\left(x_{1}\right)\right)=\sum_{n=1}^{\infty} h_{n}(0)=\sum_{n=1}^{\infty} c_{n} \leq 1 .
$$

On the other hand, if $\left(x_{1}, x_{2}\right) \in G(\bar{f}) \backslash U$, then $0=\left\|x_{2}-\bar{f}\left(x_{1}\right)\right\|<a_{n}+\varepsilon_{n}$ and $\lambda_{n}\left(x_{1}\right)=1$ for all $n \in \mathbb{N}$, so $\tilde{\psi}$ equals again the constant $d^{*}$. In fact $d^{*}$ is the supremum of $\psi$ and of $\tilde{\psi}$, and is easily seen to be attained in $G(\bar{f})$ and $G(\bar{f}) \backslash U$, respectively.
To show that $\psi$ and $\tilde{\psi}$ are equal outside $U$ take $\left(x_{1}, x_{2}\right) \in E$. By a remark after the definition of $\psi$ and $\tilde{\psi}$, we can assume $x_{1} \in U_{0}$.

Claim 2.6. For every $n \in \mathbb{N}$, if $x_{1} \in V_{n} \backslash V_{n-1}$ and if $x_{2} \in E_{2}$ is such that $\left\|x_{2}-\bar{f}\left(x_{1}\right)\right\| \geq b_{n}+\varepsilon_{n}$, then $\psi\left(x_{1}, x_{2}\right)=\tilde{\psi}\left(x_{1}, x_{2}\right)$.
Proof of Claim. If $\left\|x_{2}-\bar{f}\left(x_{1}\right)\right\| \geq b_{n}+\varepsilon_{n}$ we have that $\psi_{k}\left(x_{1}, x_{2}\right)=0$ for all $k \geq n$. So we have to see that $\lambda_{k}\left(x_{1}\right)=1$ for $k=1, \ldots, n-1$. But this is clear since $x \notin V_{n-1}$ and hence $x \notin V_{k}$ for any $k=1, \ldots, n-1$.

Now, we can conclude that for each $n \in \mathbb{N}, \psi=\tilde{\psi}$ on the set

$$
\left(\left(V_{n} \backslash V_{n-1}\right) \times E_{2}\right) \backslash U_{n} \supseteq\left(\left(V_{n} \backslash V_{n-1}\right) \times E_{2}\right) \backslash U
$$

Since $\bigcup_{n=1}^{\infty} V_{n} \backslash V_{n-1}=U_{0}$ and $\bigcup_{n=1}^{\infty} U_{n} \subseteq U$, it follows that $\psi$ is equal to $\tilde{\psi}$ outside $U$. Finally, to obtain functions $\phi$ and $\tilde{\phi}$ with the desired properties it is sufficient to set

$$
\begin{aligned}
& \phi\left(x_{1}, x_{2}\right)=\psi\left(x_{1}, x_{2}\right)+1-d^{*} \\
& \tilde{\phi}\left(x_{1}, x_{2}\right)=\tilde{\psi}\left(x_{1}, x_{2}\right)+1-d^{*}
\end{aligned}
$$

for all $\left(x_{1}, x_{2}\right) \in E$. This ensures that the supremum, which is attained precisely on $G(\bar{f})$ for $\phi$ and precisely on $G(\bar{f}) \backslash U$ for $\tilde{\phi}$, is equal to 1 .

In the proof of Theorem [2.1, we will employ a well known fact stating that the identity mapping perturbed by a contracting mapping is a homeomorphism (even a diffeomorphism provided that the contracting mapping is smooth). This fact is stated and proved in Lemma 3 of Renz's Ph.D. thesis [58.

Lemma 2.7. Let $E_{1}$ be a normed linear space and $E_{2}$ be a Banach space. Let $E=E_{1} \times E_{2}$ and let $d: E \rightarrow E_{2}$ be a continuous mapping satisfying the following condition

$$
\left\|d\left(x_{1}, x_{2}\right)-d\left(x_{1}, x_{2}^{\prime}\right)\right\| \leq \frac{1}{2}\left\|x_{2}-x_{2}^{\prime}\right\|
$$

for all $x_{1} \in E_{1}$ and $x_{2}, x_{2}^{\prime} \in E_{2}$. Then the mapping defined by $h\left(x_{1}, x_{2}\right)=\left(x_{1}, x_{2}-d\left(x_{1}, x_{2}\right)\right)$ is a homeomorphism of $E$ onto itself. Moreover, $h$ is a $C^{p}$ diffeomorphism when restricted to any open set (onto its image) on which $d$ is $C^{p}$ smooth.

Let us now present the proof of Theorem 2.1.
Proof of Theorem 2.1. Basically, we will follow the proof of Theorem 1 of Renz's Ph.D. thesis [58]. First, we apply Lemma 2.2 to $f: X_{1} \rightarrow E_{2}$ to obtain a continuous mapping $\bar{f}: E_{1} \rightarrow E_{2}$ such that $\bar{f} \mid X_{1}=f$ and $\bar{f} \mid E_{1} \backslash X_{1}$ is $C^{p}$ smooth. Then, we apply Lemma 2.3 to the mapping $\bar{f}$ to obtain a mapping $F: \mathbb{R} \times E_{1} \rightarrow E_{2}$ satisfying conditions (1)-(4) of Lemma 2.3. Next, we apply Lemma 2.5 to $\bar{f}$ to obtain functions $\phi$ and $\tilde{\phi}$ satisfying conditions (1)-(4) of Lemma 2.5. Now, we define

$$
\begin{aligned}
& d\left(x_{1}, x_{2}\right)=F\left(\phi\left(x_{1}, x_{2}\right), x_{1}\right) \\
& \tilde{d}\left(x_{1}, x_{2}\right)=F\left(\tilde{\phi}\left(x_{1}, x_{2}\right), x_{1}\right) .
\end{aligned}
$$

Let us check that Lemma 2.7 is applicable to $d$ and $\tilde{d}$ so that

$$
h\left(x_{1}, x_{2}\right)=\left(x_{1}, x_{2}-d\left(x_{1}, x_{2}\right)\right)
$$

and

$$
\varphi\left(x_{1}, x_{2}\right)=\left(x_{1}, x_{2}-\tilde{d}\left(x_{1}, x_{2}\right)\right)
$$

are homeomorphisms (which, additionally, will satisfy the conditions enumerated in Theorem 2.1).
Both functions $d$ and $\tilde{d}$ are continuous as compositions of continuous functions. We compute $D_{2} d$ and $D_{2} \tilde{d}$ to obtain

$$
\begin{aligned}
& D_{2} d\left(x_{1}, x_{2}\right)=D_{1} F\left(\phi\left(x_{1}, x_{2}\right), x_{1}\right) \circ D_{2} \phi\left(x_{1}, x_{2}\right) \\
& D_{2} \tilde{d}\left(x_{1}, x_{2}\right)=D_{1} F\left(\tilde{\phi}\left(x_{1}, x_{2}\right), x_{1}\right) \circ D_{2} \tilde{\phi}\left(x_{1}, x_{2}\right) .
\end{aligned}
$$

The estimates of the norms of $D_{1} F, D_{2} \phi$, and $D_{2} \tilde{\phi}$ yields $\left\|D_{2} d\left(x_{1}, x_{2}\right)\right\|,\left\|D_{2} \tilde{d}\left(x_{1}, x_{2}\right)\right\| \leq \frac{1}{4}$ when $\left(x_{1}, x_{2}\right) \notin G(\bar{f})$. Since $d$ and $\tilde{d}$ are continuous and $E \backslash G(\bar{f})$ is dense in $E$, by the mean value theorem, we can write

$$
\begin{aligned}
& \left\|d\left(x_{1}, x_{2}\right)-d\left(x_{1}, x_{2}^{\prime}\right)\right\| \leq \frac{1}{4}\left\|x_{2}-x_{2}^{\prime}\right\| \\
& \left\|\tilde{d}\left(x_{1}, x_{2}\right)-\tilde{d}\left(x_{1}, x_{2}^{\prime}\right)\right\| \leq \frac{1}{4}\left\|x_{2}-x_{2}^{\prime}\right\|
\end{aligned}
$$

for all $x_{1} \in E_{1}$ and all $x_{2}, x_{2}^{\prime} \in E_{2}$. Hence, Lemma 2.7 applies and yields that $h$ and $\varphi$ are homeomorphisms.
Let us show conditions (1)-(7) of Theorem 2.1,
First, we will verify condition (1). If $\left(x_{1}, x_{2}\right) \in G(f)$, then $d\left(x_{1}, x_{2}\right)=F\left(1, x_{1}\right)=f\left(x_{1}\right)=x_{2}$ and

$$
\begin{equation*}
h\left(x_{1}, x_{2}\right)=\left(x_{1}, x_{2}-d\left(x_{1}, x_{2}\right)\right)=\left(x_{1}, 0\right) \tag{2.3}
\end{equation*}
$$

If $\left(x_{1}, x_{2}\right) \in G(f) \backslash U$, then $\tilde{d}\left(x_{1}, x_{2}\right)=F\left(1, x_{1}\right)=x_{2}$ and

$$
\begin{equation*}
\varphi\left(x_{1}, x_{2}\right)=\left(x_{1}, x_{2}-\tilde{d}\left(x_{1}, x_{2}\right)\right)=\left(x_{1}, 0\right) \tag{2.4}
\end{equation*}
$$

Condition (3) is obvious. Since $\phi$ and $\tilde{\phi}$ are equal outside $U$ we obtain (2).
Let us see that $d \mid E \backslash G(f)$ and $\tilde{d} \mid E \backslash(G(f) \backslash U)$ are $C^{p}$ diffeomorphisms. If $\left(x_{1}, x_{2}\right) \notin G(\bar{f})$ then $\phi$ and $\tilde{\phi}$ are $C^{p}$ smooth and $\phi\left(x_{1}, x_{2}\right), \tilde{\phi}\left(x_{1}, x_{2}\right)<1$ by condition (2) and (1) of Lemma 2.5, It follows that $F\left(\phi\left(x_{1}, x_{2}\right), x_{1}\right)$ and $F\left(\tilde{\phi}\left(x_{1}, x_{2}\right), x_{1}\right)$ are $C^{p}$ smooth in a neighborhood of $\left(x_{1}, x_{2}\right)$. Thus $h$ and $\varphi$ are $C^{p}$ on $E \backslash G(\bar{f})$.
On the other hand, we have $\phi \mid G(\bar{f}) \backslash G(f)=1$. Then, by continuity of $\phi$ and condition (1) of Lemma [2.3, we infer that $F\left(\phi\left(x_{1}, x_{2}\right), x_{1}\right)=\bar{f}\left(x_{1}\right)$ and, consequently, $h\left(x_{1}, x_{2}\right)=\left(x_{1}, x_{2}-\bar{f}\left(x_{1}\right)\right)$ in a neighborhood of $G(\bar{f}) \backslash G(f)$. We have proved that $d$ is $C^{p}$ smooth on $E \backslash G(f)$.
It remains to show that $\varphi \mid U$ is $C^{p}$ smooth. By condition (1) of Lemma 2.5, we have $\tilde{\phi} \mid U<1$; by condition (2) of Lemma [2.3, $\tilde{d} \mid U$ is $C^{p}$ smooth. The proof that $\tilde{d}$ and, therefore, $\varphi$ restricted to $E \backslash(G(f) \backslash U)$ is $C^{p}$ smooth is complete.
Now, Lemma 2.7 tells us that $h$ and $\varphi$ are $C^{p}$ diffeomorphisms of $E \backslash G(f)$ onto $h(E \backslash G(f))$ and $E \backslash(G(f) \backslash U)$ onto $\varphi(E \backslash(G(f) \backslash U))$. So to get (4) and (5) it is sufficient to show that $h(G(f))=X_{1} \times\{0\}$ and $\varphi(G(f) \backslash U)=\left(X_{1} \times\{0\}\right) \backslash h(U)$. The first equality is clear from equation (2.3). For the second one, observe that (2.4) tells us that $\varphi(G(f) \backslash U)=\left(X_{1} \backslash U_{0}\right) \times\{0\}=Y_{1} \times\{0\}$. So, we must check that

$$
\left(X_{1} \times\{0\}\right) \backslash h(U)=\left(X_{1} \backslash U_{0}\right) \times\{0\}
$$

or, what is the same, that $\pi_{1}(h(U))=U_{0}$. The latter follows from condition (3) and the fact that $\pi_{1}(U)=U_{0}$.
Let us finish the proof by showing (6) and (7). Firstly let us check that $\left\|h^{-1}\left(x_{1}, x_{2}\right)-\varphi^{-1}\left(x_{1}, x_{2}\right)\right\| \leq \varepsilon$ for every $\left(x_{1}, x_{2}\right) \in E_{1} \times E_{2}$. Since $h^{-1}$ preserves the first coordinate, we can write $h^{-1}\left(x_{1}, x_{2}\right)=$ $\left(x_{1}, y_{2}\right)$ and $\varphi^{-1}\left(x_{1}, x_{2}\right)=\left(x_{1}, z_{2}\right)$ where $y_{2}, z_{2} \in E_{2}$ are such that

$$
\begin{aligned}
& y_{2}-d\left(x_{1}, y_{2}\right)=x_{2} \\
& z_{2}-\tilde{d}\left(x_{1}, z_{2}\right)=x_{2} .
\end{aligned}
$$

We then have that $\left\|h^{-1}\left(x_{1}, x_{2}\right)-\varphi^{-1}\left(x_{1}, x_{2}\right)\right\| \leq \varepsilon$ if and only if $\left\|y_{2}-z_{2}\right\| \leq \varepsilon$ and if and only if $\left\|d\left(x_{1}, y_{2}\right)-\tilde{d}\left(x_{1}, z_{2}\right)\right\| \leq \varepsilon$. Since $r \mapsto F\left(r, x_{1}\right)$ is 1 -Lipschitz, this is true if $\left|\phi\left(x_{1}, y_{2}\right)-\tilde{\phi}\left(x_{1}, z_{2}\right)\right| \leq \epsilon$,
or what is the same if $\left|\psi\left(x_{1}, y_{2}\right)-\tilde{\psi}\left(x_{1}, z_{2}\right)\right| \leq \epsilon$. And this is the case because

$$
\begin{aligned}
\left|\psi\left(x_{1}, y_{2}\right)-\tilde{\psi}\left(x_{1}, z_{2}\right)\right| & \leq\left|\psi\left(x_{1}, y_{2}\right)\right|+\left|\tilde{\psi}\left(x_{1}, z_{2}\right)\right|=\left|\sum_{n=1}^{\infty} \psi_{n}\left(x_{1}, y_{2}\right)\right|+\left|\sum_{n=1}^{\infty} \lambda_{n}\left(x_{1}\right) \psi_{n}\left(x_{1}, z_{2}\right)\right| \leq \\
& \leq \sum_{n=1}^{\infty} c_{n}+\sum_{n=1}^{\infty} c_{n} \leq \frac{\epsilon}{2}+\frac{\epsilon}{2}=\epsilon
\end{aligned}
$$

Secondly, let us see that for every $\eta>0$ there exists $\delta>0$ such that if $\left\|\left(x_{1}, x_{2}\right)-\left(x_{1}, x_{2}^{\prime}\right)\right\|=$ $\left\|x_{2}-x_{2}^{\prime}\right\| \leq \delta$ then $\left\|h^{-1}\left(x_{1}, x_{2}\right)-h^{-1}\left(x_{1}, x_{2}^{\prime}\right)\right\| \leq \eta$. It will be enough to set $\delta=\frac{\eta}{2}$. Indeed, take $\left(x_{1}, x_{2}\right),\left(x_{1}, x_{2}^{\prime}\right) \in E_{1} \times E_{2}$ such that $\left\|\left(x_{1}, x_{2}\right)-\left(x_{1}, x_{2}^{\prime}\right)\right\|=\left\|x_{2}-x_{2}^{\prime}\right\| \leq \frac{\eta}{2}$. Write $h^{-1}\left(x_{1}, x_{2}\right)=\left(x_{1}, y_{2}\right)$ and $h^{-1}\left(x_{1}, x_{2}^{\prime}\right)=\left(x_{1}, y_{2}^{\prime}\right)$ where $y_{2}, y_{2}^{\prime} \in E_{2}$ are such that

$$
\begin{aligned}
& y_{2}-d\left(x_{1}, y_{2}\right)=x_{2} \\
& y_{2}^{\prime}-\tilde{d}\left(x_{1}, y_{2}^{\prime}\right)=x_{2}^{\prime} .
\end{aligned}
$$

Then we have that

$$
\begin{aligned}
\left\|h^{-1}\left(x_{1}, x_{2}\right)-h^{-1}\left(x_{1}, x_{2}^{\prime}\right)\right\| & =\left\|y_{2}-y_{2}^{\prime}\right\| \leq\left\|x_{2}-x_{2}^{\prime}\right\|+\left\|d\left(x_{1}, y_{2}\right)-d\left(x_{1}, y_{2}^{\prime}\right)\right\| \leq \\
& \leq\left\|x_{2}-x_{2}^{\prime}\right\|+\frac{1}{2}\left\|y_{2}-y_{2}^{\prime}\right\|,
\end{aligned}
$$

which implies that $\left\|y_{2}-y_{2}^{\prime}\right\| \leq 2\left\|x_{2}-x_{2}^{\prime}\right\|=\eta$, and the proof is complete.

### 2.2. An extracting scheme tailored for closed subsets of a subspace of infinite codimension.

In this subsection we will establish the following.
Theorem 2.8. Let $E_{1}$ and $E_{2}$ be Banach spaces such that $E_{2}$ is infinite-dimensional and admits a (not necessarily equivalent) $C^{p}$ smooth norm, where $p \in \mathbb{N} \cup\{\infty\}$. Define $E=E_{1} \times E_{2}$ and, for $i=1,2$, write $\pi_{i}: E \rightarrow E_{i}$ for the natural projections, that is, $\pi_{i}\left(x_{1}, x_{2}\right)=x_{i}$ for $\left(x_{1}, x_{2}\right) \in E$. Let $W_{1}$ be an open subset of $E_{1}$, and $\psi: W_{1} \rightarrow[0, \infty)$ be a continuous function such that $\psi$ is of class $C^{p}$ on $\psi^{-1}(0, \infty)$. Denote $K=\psi^{-1}(0) \times\{0\}$. Then, there exists a $C^{p}$ diffeomorphism $h$ from $\left(W_{1} \times E_{2}\right) \backslash K$ onto $W_{1} \times E_{2}$ which satisfies $\pi_{1} \circ h=h$ and is the identity off of a certain open subset $U$ of $W_{1} \times E_{2}$. Specifically, the set $U$ is defined as follows

$$
U:=\left\{x=\left(x_{1}, x_{2}\right) \in W_{1} \times E_{2}: S\left(\psi\left(x_{1}\right), \omega\left(x_{2}\right)\right)<1\right\}
$$

where $S$ is a certain $C^{\infty}$ norm on $\mathbb{R}^{2}$ and $\omega: E_{2} \rightarrow[0, \infty)$ is a certain (not necessarily symmetric) subadditive and positive-homogeneous functional of class $C^{p}$ on $E_{2} \backslash\{0\}$; see Lemmas 2.10 and 2.13 for precise definitions.

By a $C^{p}$ smooth norm on $E_{2}$ we mean a (possibly nonequivalent) norm on $E_{2}$ which is of class $C^{p}$ on $E_{2} \backslash\{0\}$.
We will need to use the following three auxiliary results from (4), [5].
Lemma 2.9. Let $F:(0, \infty) \longrightarrow[0, \infty)$ be a continuous function such that, for every $\beta \geq \alpha>0$,

$$
F(\beta)-F(\alpha) \leq \frac{1}{2}(\beta-\alpha), \quad \text { and } \quad \limsup _{t \rightarrow 0^{+}} F(t)>0
$$

Then there exists a unique $\alpha>0$ such that $F(\alpha)=\alpha$.
Proof. See [4, Lemma 2].

It is not known whether every infinite-dimensional Banach space with a $C^{1}$ equivalent norm possesses a $C^{1}$ smooth non-complete norm $\sqrt[1]{1}$ The following lemma shows that for every Banach space with a $C^{p}$ smooth norm there exists a kind of $C^{p}$ asymmetric non-complete subadditive functional which successfully replaces the smooth non-complete norm in Bessaga's technique [14] for extracting points.

Lemma 2.10. Let $\left(E_{2},\|\cdot\|\right)$ be an infinite-dimensional Banach space which admits a (not necessarily equivalent) $C^{p}$ smooth norm, where $p \in \mathbb{N} \cup\{\infty\}$. Then there exists a continuous function $\omega: E_{2} \longrightarrow$ $[0, \infty)$ which is $C^{p}$ smooth on $E_{2} \backslash\{0\}$ and satisfies the following properties:
(1) $\omega(x+y) \leq \omega(x)+\omega(y)$, and, consequently, $\omega(x)-\omega(y) \leq \omega(x-y)$, for every $x, y \in E_{2}$;
(2) $\omega(r x)=r \omega(x)$ for every $x \in E_{2}$, and $r \geq 0$;
(3) $\omega(x)=0$ if and only if $x=0$;
(4) $\omega\left(\sum_{k=1}^{\infty} z_{k}\right) \leq \sum_{k=1}^{\infty} \omega\left(z_{k}\right)$ for every convergent series $\sum_{k=1}^{\infty} z_{k}$ in $\left(E_{2},\|\cdot\|\right)$; and
(5) for every $\varepsilon>0$, there exists a sequence of vectors $\left(y_{k}\right) \subset E_{2}$ such that

$$
\omega\left(y_{k}\right) \leq \frac{\varepsilon}{4^{k+1}},
$$

for every $k \in \mathbb{N}$, and

$$
\liminf _{n \rightarrow \infty} \omega\left(y-\sum_{j=1}^{n} y_{j}\right)>0
$$

for every $y \in E_{2}$.
Notice that $\omega$ need not be a norm in $E_{2}$, as in general we have $\omega(x) \neq \omega(-x)$.
Proof. See [5, Lemma 2.3]
Using the properties of the functional $\omega$ we can construct an extracting curve as follows.
Lemma 2.11. Let $\left(E_{2},\|\cdot\|\right)$ be a Banach space, and let $\omega$ be a functional satisfying conditions (1), (2), and (5) of Lemma 2.10. Then there exists a $C^{\infty}$ curve $\gamma:(0, \infty) \longrightarrow E_{2}$ such that
(1) $\omega(\gamma(\alpha)-\gamma(\beta)) \leq \frac{1}{2}(\beta-\alpha)$ if $\beta \geq \alpha>0$;
(2) $\lim \sup _{t \rightarrow 0^{+}} \omega(y-\gamma(t))>0$ for every $y \in E_{2}$; and
(3) $\gamma(t)=0$ if $t \geq 1$.

Proof. Let $\theta:[0, \infty) \longrightarrow[0,1]$ be a non-increasing $C^{\infty}$ function such that $\theta=1$ on $[0,1 / 2], \theta=0$ on $[1, \infty)$ and $\sup \left\{\left|\theta^{\prime}(t)\right|: t \in[0, \infty)\right\} \leq 4$. Let us choose a sequence of vectors $\left(y_{k}\right) \subset E_{2}$ which satisfies condition (5) of Lemma 2.10 for $\varepsilon=1$, and define $\gamma:(0, \infty) \longrightarrow E_{2}$ by the following formula

$$
\gamma(t)=\sum_{k=1}^{\infty} \theta\left(2^{k-1} t\right) y_{k}
$$

It is not difficult to check that this curve satisfies the properties of the statement. See [5, Lemma 2.5] for details.

We will also need a technical tool (see, for instance, [8, Lemmas 2.27 and 2.28]) that allows us to obtain, on the product space $E_{1} \times E_{2}$, a norm which preserves the smoothness properties that the corresponding norms of the factors may have. Notice that the natural formula $\left(\left\|x_{1}\right\|_{1}^{2}+\left\|x_{2}\right\|_{2}^{2}\right)^{1 / 2}$ defines a $C^{1}$ norm in $E_{1} \times E_{2} \backslash\{0\}$ but, in general, this norm will not be $C^{2}$ on this set, even if $\|\cdot\|_{1}$ and $\|\cdot\|_{2}$ are $C^{\infty}$ on $E_{1} \backslash\{0\}$ and $E_{2} \backslash\{0\}$, respectively, because the function $x_{2} \mapsto\left\|x_{2}\right\|_{2}^{2}$ may not be $C^{2}$ smooth on all of $E_{2}$ even though it is $C^{\infty}$ smooth on $E_{2} \backslash\{0\}$. As a matter of fact, it is not difficult to show that, for every Banach space $(E,\|\cdot\|)$, if $\|\cdot\|^{2}$ is twice Fréchet differentiable at 0 , then $E$ is isomorphic to a Hilbert space; see, for instance [33, Exercise 10.4, pp. 475-476].

[^1]Definition 2.12. We will say that a subset $\mathcal{S}$ of the plane $\mathbb{R}^{2}$ is a smooth square provided that:
(i) $\mathcal{S} \subset \mathbb{R}^{2}$ is a bounded, symmetric convex body with $0 \in \operatorname{int}(\mathcal{S})$, and whose boundary $\partial \mathcal{S}$ is $C^{\infty}$ smooth.
(ii) $(x, y) \in \partial \mathcal{S} \Leftrightarrow\left(\epsilon_{1} x, \epsilon_{2} y\right) \in \partial \mathcal{S}$ for each couple $\left(\epsilon_{1}, \epsilon_{2}\right) \in\{-1,1\}^{2}$ (that is, $\mathcal{S}$ is symmetric about the cordinate axes).
(iii) $\left[-\frac{1}{2}, \frac{1}{2}\right] \times\{-1,1\} \cup\{-1,1\} \times\left[-\frac{1}{2}, \frac{1}{2}\right] \subset \partial \mathcal{S}$.
(iv) $\mathcal{S} \subset[-1,1] \times[-1,1]$.

Of course, it is elementary to produce smooth squares in $\mathbb{R}^{2}$.
The following lemma enumerates the essential properties of a smooth square. Recall that the Minkowski functional of a convex body $\mathcal{A}$ such that $0 \in \operatorname{int}(\mathcal{A})$ is defined by

$$
\mu_{\mathcal{A}}(x)=\inf \left\{t>0: \frac{1}{t} x \in \mathcal{A}\right\} .
$$

Lemma 2.13. Let $\mathcal{S} \subset \mathbb{R}^{2}$ be a smooth square. Then its Minkowski functional $\mu_{\mathcal{S}}: \mathbb{R}^{2} \rightarrow \mathbb{R}$ is a $C^{\infty}$ smooth norm on $\mathbb{R}^{2}$ such that, for every $(x, y) \in \mathbb{R}^{2}$, we have
(1) $\mu_{\mathcal{S}}(x, y) \leq|x|+|y| \leq 2 \mu_{\mathcal{S}}(x, y)$;
(2) $\max (|x|,|y|) \leq \mu_{\mathcal{S}}(x, y) \leq 2 \max (|x|,|y|)$;
(3) $\mu_{\mathcal{S}}(0, y)=|y|, \quad \mu_{\mathcal{S}}(x, 0)=|x|$;
(4) For every $\left(x_{0}, y_{0}\right) \in \mathbb{R}^{2} \backslash\{(0,0)\}$, there exists $\sigma>0$ so that

$$
\mu_{\mathcal{S}}(x, y)=|x| \quad \text { if } \quad \max \left(\left|x-x_{0}\right|,|y|\right) \leq \sigma \quad \text { and } \quad \mu_{\mathcal{S}}(x, y)=|y| \quad \text { if } \quad \max \left(|x|,\left|y-y_{0}\right|\right) \leq \sigma .
$$

(5) The functions $(0, \infty) \ni t \rightarrow \mu_{\mathcal{S}}(x, t y)$ and $(0, \infty) \ni t \rightarrow \mu_{\mathcal{S}}(t x, y)$ are both nondecreasing.

Note that property (4) (which is related to properties (iii) and (iv) of Definition (2.12) means that every sphere of $\left(\mathbb{R}^{2}, \mu_{\mathcal{S}}\right)$ centered at the origin, which coincides with $\lambda(\partial \mathcal{S})$ for some $\lambda>0$, is orthogonal to the coordinate axes and is locally flat on a neighborhood of the intersection of $\lambda(\partial \mathcal{S})$ with the lines $\{x=0\} \cup\{y=0\}$. By using this property it is easy to show that, for any couple of Banach spaces $\left(E_{1},\|\cdot\|_{1}\right)$ and $\left(E_{2},\|\cdot\|_{2}\right)$ with $C^{p}$ smooth norms, the expression

$$
\mu_{\mathcal{S}}\left(\left\|x_{1}\right\|_{1},\left\|x_{2}\right\|_{2}\right)
$$

defines an equivalent norm of class $C^{p}$ in $E_{1} \times E_{2}$.
Proof. Properties (1) - (3) and (5) are easy to show. Let us prove (4). Assume for instance that $x_{0} \neq 0$, and set $\sigma=\left|x_{0}\right| / 4$. If $\max \left(\left|x-x_{0}\right|,|y|\right) \leq \sigma$, then we have

$$
\frac{|y|}{|x|} \leq \frac{\sigma}{\left|x_{0}\right|-\sigma}=\frac{\left|x_{0}\right| / 4}{\left|x_{0}\right|-\left|x_{0}\right| / 4}=\frac{1}{3}<\frac{1}{2},
$$

hence

$$
\left(\frac{x}{|x|}, \frac{y}{|x|}\right) \in\{-1,1\} \times\left[-\frac{1}{2}, \frac{1}{2}\right] \subset \partial \mathcal{S},
$$

and it follows that $\mu_{\mathcal{S}}(x, y)=|x|$.
The lemma below shows how, with the help of a smooth square, we can combine the given $C^{p}$ smooth function $\psi: W_{1} \rightarrow[0, \infty)$ together with the $C^{p}$ smooth functional $\omega: E_{2} \rightarrow[0, \infty)$ obtained in Lemma 2.10, in order to obtain a $C^{p}$ smooth function on $W_{1} \times E_{2}$ which behaves more or less like $\psi\left(x_{1}\right)+\omega\left(x_{2}\right)$ (or, equivalently, like $\left(\psi\left(x_{1}\right)^{2}+\omega\left(x_{2}\right)^{2}\right)^{1 / 2}$ ).

Lemma 2.14. Let $E=E_{1} \times E_{2}$ be a Banach space, $\rho_{1}: E_{1} \rightarrow[0, \infty)$ and $\rho_{2}: E_{2} \rightarrow[0, \infty)$ continuous functions which are of class $C^{p}$ on $E_{1} \backslash \rho_{1}^{-1}(0)$ and $E_{2} \backslash \rho_{1}^{-1}(0)$, respectively. Then, for any smooth square $\mathcal{S}$ of $\mathbb{R}^{2}$, the function $\rho: E_{1} \times E_{2} \rightarrow[0, \infty)$ defined by

$$
\rho(x)=\rho\left(x_{1}, x_{2}\right)=\mu_{\mathcal{S}}\left(\rho_{1}\left(x_{1}\right), \rho_{2}\left(x_{2}\right)\right), \quad x=\left(x_{1}, x_{2}\right) \in E_{1} \times E_{2},
$$

is continuous on $E$ and of class $C^{p}$ on $E \backslash\left(\rho_{1}^{-1}(0) \times \rho_{2}^{-1}(0)\right)$.
The same is true if we replace $E_{1}$ with an open subset $W_{1}$ of $E_{1}$.
Proof. It is clear that $\rho$ is continuous on $E$, and that it is $C^{p}$ smooth on $\left\{\left(x_{1}, x_{2}\right) \in E: \rho_{1}\left(x_{1}\right) \neq 0 \neq\right.$ $\left.\rho_{2}\left(x_{2}\right)\right\}$. Let us see that $\rho$ is also $C^{p}$ smooth on a neighborhood of the set

$$
\left(\left\{\left(x_{1}, x_{2}\right) \in E: \rho_{1}\left(x_{1}\right)=0\right\} \cup\left\{\left(x_{1}, x_{2}\right) \in E: \rho_{2}\left(x_{2}\right)=0\right\}\right) \backslash\left(\rho_{1}^{-1}(0) \times \rho_{2}^{-1}(0)\right) .
$$

Suppose for instance that $\rho_{1}\left(x_{1}\right) \neq 0=\rho_{2}\left(x_{2}\right)$. Then, by continuity of $\rho_{1}, \rho_{2}$ and by property (4) of Lemma 2.13, there exist a neighborhood $U$ of the point $\left(x_{1}, x_{2}\right)$ such that $U \subset\left\{\left(y_{1}, y_{2}\right) \in E: \rho_{1}\left(y_{1}\right) \neq\right.$ $0\}$ and

$$
\rho\left(y_{1}, y_{2}\right)=\rho_{1}\left(y_{1}\right)
$$

for all $\left(y_{1}, y_{2}\right) \in U$. It follows that $\rho$ is of class $C^{p}$ on $U$. The case $\rho_{1}\left(x_{1}\right)=0 \neq \rho_{2}\left(x_{2}\right)$ can be treated similarly.
Now we are ready to prove Theorem 2.8,
Proof of Theorem 2.8. From now on we will fix a smooth square $\mathcal{S}$ on $\mathbb{R}^{2}$, and we will denote

$$
S=\mu_{\mathcal{S}} .
$$

Thus, by Lemma 2.14 applied to $\rho_{1}=\psi$ and $\rho_{2}=\omega$ (recall that $\omega$ was constructed in Lemma 2.10), the function

$$
\rho\left(x_{1}, x_{2}\right):=S\left(\psi\left(x_{1}\right), \omega\left(x_{2}\right)\right)
$$

is continuous on $W_{1} \times E_{2}$ and of class $C^{p}$ on $\left(W_{1} \times E_{2}\right) \backslash\left(\psi^{-1}(0) \times\{0\}\right)$.
Let us define $h:\left(W_{1} \times E_{2}\right) \backslash K \rightarrow E$ by

$$
h\left(x_{1}, x_{2}\right)=\left(x_{1}, x_{2}+\gamma \circ \rho\left(x_{1}, x_{2}\right)\right)=\left(x_{1}, x_{2}+\gamma\left(S\left(\psi\left(x_{1}\right), \omega\left(x_{2}\right)\right)\right)\right), \quad\left(x_{1}, x_{2}\right) \in\left(W_{1} \times E_{2}\right) \backslash K,
$$

where $\gamma$ is provided by Lemma 2.11. Note that

$$
K=\psi^{-1}(0) \times\{0\}=\rho^{-1}(0)
$$

Let $\left(y_{1}, y_{2}\right)$ be an arbitrary point of $W_{1} \times E_{2}$, and let $F_{y_{1}, y_{2}}:(0, \infty) \longrightarrow[0, \infty)$ be defined by

$$
\begin{equation*}
F_{y_{1}, y_{2}}(\alpha)=\rho\left(y_{1}, y_{2}-\gamma(\alpha)\right)=S\left(\psi\left(y_{1}\right), \omega\left(y_{2}-\gamma(\alpha)\right)\right) \tag{2.5}
\end{equation*}
$$

for $\alpha>0$. Let us see that $F_{y_{1}, y_{2}}(\alpha)$ satisfies the conditions of Lemma 2.9, As for the first condition, we consider two cases: if $\omega\left(y_{2}-\gamma(\beta)\right) \leq \omega\left(y_{2}-\gamma(\alpha)\right)$ then, since the function $(0, \infty) \ni t \mapsto S\left(\psi\left(y_{1}\right), t\right)$ is increasing (see condition (5) of Lemma 2.13), we have that

$$
S\left(\psi\left(y_{1}\right), \omega\left(y_{2}-\gamma(\beta)\right)\right) \leq S\left(\psi\left(y_{1}\right), \omega\left(y_{2}-\gamma(\alpha)\right)\right)
$$

and therefore

$$
F_{y_{1}, y_{2}}(\beta)-F_{y_{1}, y_{2}}(\alpha) \leq 0 \leq \frac{1}{2}(\beta-\alpha)
$$

trivially for all $\beta \geq \alpha>0$. Otherwise, we have $\omega\left(y_{2}-\gamma(\beta)\right)-\omega\left(y_{2}-\gamma(\alpha)\right)>0$, and therefore, using the fact that $S$ is a norm in $\mathbb{R}^{2}$, condition (3) of Lemma 2.13, the properties of the functional $\omega$, and condition (1) of Lemma 2.11, we obtain

$$
\begin{aligned}
& S\left(\psi\left(y_{1}\right), \omega\left(y_{2}-\gamma(\beta)\right)\right)-S\left(\psi\left(y_{1}\right), \omega\left(y_{2}-\gamma(\alpha)\right)\right) \leq S\left(0, \omega\left(y_{2}-\gamma(\beta)\right)-\omega\left(y_{2}-\gamma(\alpha)\right)\right)= \\
& \omega\left(y_{2}-\gamma(\beta)\right)-\omega\left(y_{2}-\gamma(\alpha)\right) \leq \omega\left(y_{2}-\gamma(\beta)-\left(y_{2}-\gamma(\alpha)\right)\right)=\omega(\gamma(\alpha)-\gamma(\beta)) \leq \frac{1}{2}(\beta-\alpha)
\end{aligned}
$$

for every $\beta \geq \alpha>0$. In either case we have that

$$
\begin{equation*}
F_{y_{1}, y_{2}}(\beta)-F_{y_{1}, y_{2}}(\alpha) \leq \frac{1}{2}(\beta-\alpha) \tag{2.6}
\end{equation*}
$$

for all $\beta \geq \alpha>0$.
On the other hand, by condition (2) of Lemma 2.11 we know that

$$
\limsup _{\alpha \rightarrow 0^{+}} \omega\left(y_{2}-\gamma(\alpha)\right)>0,
$$

and therefore, by condition (2) of Lemma 2.13, we have

$$
\limsup _{\alpha \rightarrow 0^{+}} F_{y_{1}, y_{2}}(\alpha)=\limsup _{\alpha \rightarrow 0^{+}} S\left(\psi\left(y_{1}\right), \omega\left(y_{2}-\gamma(\alpha)\right)\right) \geq \limsup _{\alpha \rightarrow 0^{+}} \omega\left(y_{2}-\gamma(\alpha)\right)>0,
$$

so that $F_{y_{1}, y_{2}}$ also satisfies the second condition of Lemma 2.9 ,
Then, applying Lemma 2.9, we deduce that the equation $F_{y_{1}, y_{2}}(\alpha)=\alpha$ has a unique solution. This means that, for each $\left(y_{1}, y_{2}\right) \in W_{1} \times E_{2}$, a number $\alpha\left(y_{1}, y_{2}\right)>0$ with the property

$$
\begin{equation*}
S\left(\psi\left(y_{1}\right), \omega\left(y_{2}-\gamma\left(\alpha\left(y_{1}, y_{2}\right)\right)\right)\right)=\rho\left(y_{1}, y_{2}-\gamma\left(\alpha\left(y_{1}, y_{2}\right)\right)\right)=\alpha\left(y_{1}, y_{2}\right) \tag{2.7}
\end{equation*}
$$

is uniquely determined.
Let us see why these facts imply that $h$ is a $C^{p}$ diffeomorphism from $W_{1} \times E_{2} \backslash K$ onto $W_{1} \times E_{2}$. Assume first that $h\left(x_{1}, x_{2}\right)=\left(y_{1}, y_{2}\right)=h\left(z_{1}, z_{2}\right)$, that is to say $x_{1}=y_{1}=z_{1}$, and

$$
\begin{equation*}
x_{2}+\gamma\left(\rho\left(y_{1}, x_{2}\right)\right)=y_{2}=z_{2}+\gamma\left(\rho\left(y_{1}, z_{2}\right)\right), \tag{2.8}
\end{equation*}
$$

or equivalently

$$
\left(y_{1}, x_{2}\right)=\left(y_{1}, y_{2}-\gamma\left(\rho\left(y_{1}, x_{2}\right)\right) \quad \text { and } \quad\left(y_{1}, z_{2}\right)=\left(y_{1}, y_{2}-\gamma\left(\rho\left(y_{1}, z_{2}\right)\right) .\right.\right.
$$

Applying $\rho$ to all sides of the above equations and using (2.5), we obtain
$\rho\left(y_{1}, x_{2}\right)=\rho\left(y_{1}, y_{2}-\gamma\left(\rho\left(y_{1}, x_{2}\right)\right)\right)=F\left(\rho\left(y_{1}, x_{2}\right)\right)$ and $\rho\left(y_{1}, z_{2}\right)=\rho\left(y_{1}, y_{2}-\gamma\left(\rho\left(y_{1}, z_{2}\right)\right)\right)=F\left(\rho\left(y_{1}, z_{2}\right)\right)$.
It follows that both $\rho\left(y_{1}, x_{2}\right)$ and $\rho\left(y_{1}, z_{2}\right)$ are fixed points of $F_{y_{1}, y_{2}}$. By the uniqueness of the fixed point, we conclude that

$$
\alpha\left(y_{1}, y_{2}\right)=\rho\left(y_{1}, x_{2}\right)=\rho\left(y_{1}, z_{2}\right) .
$$

Now applying (2.8), we have

$$
x_{2}=z_{2} \quad \text { and } \quad x_{2}=y_{2}-\gamma\left(\alpha\left(y_{1}, y_{2}\right)\right) .
$$

This shows that $h$ is one to one, and also that, given $\left(y_{1}, y_{2}\right) \in W_{1} \times E_{2}$ we have

$$
h\left(y_{1}, y_{2}-\gamma\left(\alpha\left(y_{1}, y_{2}\right)\right)\right)=\left(y_{1}, y_{2}\right) .
$$

Hence $h$ is also onto, and $h^{-1}: W_{1} \times E_{2} \rightarrow W_{1} \times E_{2} \backslash K$ is given by

$$
h^{-1}\left(y_{1}, y_{2}\right)=\left(y_{1}, y_{2}-\gamma\left(\alpha\left(y_{1}, y_{2}\right)\right)\right) .
$$

It is clear that $h$ is of class $C^{p}$. In order to see that $h^{-1}$ is $C^{p}$ as well, let us define $\Phi: W_{1} \times E_{2} \times$ $(0, \infty) \longrightarrow \mathbb{R}$ by

$$
\Phi\left(y_{1}, y_{2}, \alpha\right)=\alpha-S\left(\psi\left(y_{1}\right), \omega\left(y_{2}-\gamma(\alpha)\right)\right)=\alpha-\rho\left(y_{1}, y_{2}-\gamma(\alpha)\right) .
$$

On the one hand, according to (2.7) and the fact that $S$ is a norm in $\mathbb{R}^{2}$, we have

$$
\left(\psi\left(y_{1}\right), \omega\left(y_{2}-\gamma\left(\alpha\left(y_{1}, y_{2}\right)\right)\right) \neq(0,0)\right.
$$

for every $\left(y_{1}, y_{2}\right) \in W_{1} \times E_{2}$. Since $S$ is $C^{\infty}$ smooth away from $(0,0)$, this implies that $\Phi$ is $C^{p}$ smooth on a neighborhood of every point $\left(y_{1}, y_{2}, \alpha\left(y_{1}, y_{2}\right)\right)$ in $W_{1} \times E_{2} \times(0, \infty)$. On the other hand, we know
that $F_{y_{1}, y_{2}}(\beta)-F_{y_{1}, y_{2}}(\alpha) \leq \frac{1}{2}(\beta-\alpha)$ for $\beta \geq \alpha>0$, which implies that $F_{y_{1}, y_{2}}^{\prime}(\alpha) \leq \frac{1}{2}$ for every $\alpha$ in a neighborhood of $\alpha\left(y_{1}, y_{2}\right)$, and therefore

$$
\frac{\partial \Phi\left(y_{1}, y_{2}, \alpha\right)}{\partial \alpha}=1-F_{y_{1}, y_{2}}^{\prime}(\alpha) \geq 1-1 / 2>0
$$

Hence, by the implicit function theorem, the mapping ( $y_{1}, y_{2}$ ) $\rightarrow \alpha\left(y_{1}, y_{2}\right)$ is of class $C^{p}$ on $W_{1} \times E_{2}$, and, since $\gamma$ is $C^{p}$ smooth, so is $h^{-1}$.
Finally, it is obvious that $\pi_{1} \circ h=\pi_{1}$, and the fact that $\gamma(t)=0$ whenever $t \geq 1$ implies that $h$ is the identity off of the set $\left\{x=\left(x_{1}, x_{2}\right) \in W_{1} \times E_{2}: S\left(\psi\left(x_{1}\right), \omega\left(x_{2}\right)\right)<1\right\}$.
2.3. Extracting pieces of continuous graphs of infinite codimension. Now we will prove the following extractibility result.

Theorem 2.15. Let $E=E_{1} \times E_{2}$ be a product of Banach spaces such that $E_{1}$ admits $C^{p}$ smooth partitions of unity and $E_{2}$ admits a $C^{p}$ (not necessarily equivalent) norm. Assume that $X_{1}$ is a closed subset of $E_{1}$, that $f: E_{1} \rightarrow E_{2}$ is a continuous mapping, and that $E_{2}$ is infinite-dimensional. Define

$$
X=\left\{\left(x_{1}, x_{2}\right) \in E_{1} \times E_{2}: x_{1} \in X_{1}, x_{2}=f\left(x_{1}\right)\right\}
$$

Let $U$ be an open subset of $E$ and $\varepsilon>0$. Then there exists a $C^{p}$ diffeomorphism $g$ from $E \backslash X$ onto $E \backslash(X \backslash U)$ such that $g$ is the identity on $(E \backslash U) \backslash X$ and moves no point more than $\varepsilon$.

Proof. We may of course assume $U \cap X \neq \emptyset$ (as otherwise the result holds trivially with $g$ equal to the identity map).
Claim 2.16. It is sufficient to prove the result for $f=0$ and such that the extracting diffeomorphism preserves the first coordinate, that is $g\left(x_{1}, x_{2}\right)=\left(x_{1}, \pi_{2}\left(g\left(x_{1}, x_{2}\right)\right)\right)$.
Proof. Let $h$ and $\varphi$ be homeomorphisms given by Theorem 2.1 such that $\left\|\varphi^{-1}(x)-h^{-1}(x)\right\| \leq \frac{\varepsilon}{2}$ for every $x \in E$. By the uniform continuity of $h^{-1}\left(x_{1}, x_{2}\right)$ with respect to the second variable $x_{2} \in E_{2}$, we may choose $\delta>0$ such that if $\left\|\left(x_{1}, x_{2}\right)-\left(x_{1}, x_{2}^{\prime}\right)\right\| \leq \delta$ then

$$
\left\|h^{-1}\left(x_{1}, x_{2}\right)-h^{-1}\left(x_{1}, x_{2}^{\prime}\right)\right\| \leq \frac{\varepsilon}{2} .
$$

Assuming the result is true for $f=0$ we can find a $C^{p}$ diffeomorphism $g: E \backslash\left(X_{1} \times\{0\}\right) \rightarrow$ $E \backslash\left(\left(X_{1} \times\{0\}\right) \backslash h(U)\right)$ such that $g$ is the identity on $(E \backslash h(U)) \backslash\left(X_{1} \times\{0\}\right)$, moves no point more than $\delta$ and preserves the first coordinate. Then the composition

$$
\varphi^{-1} \circ g \circ h: E \backslash X \rightarrow(E \backslash(X \backslash U)
$$

defines a $C^{p}$ diffeomorphism with the required properties. Observe that

$$
\|g(h(x))-h(x)\|=\left\|\left(x_{1}, \pi_{2}(g(h(x)))\right)-\left(x_{1}, \pi_{2}(h(x))\right)\right\| \leq \delta
$$

hence

$$
\left\|\varphi^{-1}(g(h(x)))-x\right\| \leq\left\|\varphi^{-1}(g(h(x)))-h^{-1}(g(h(x)))\right\|+\left\|h^{-1}(g(h(x)))-x\right\| \leq \frac{\varepsilon}{2}+\frac{\varepsilon}{2}=\varepsilon
$$

for every $x \in E \backslash X$.
So it will be enough to see that if $X_{1}$ is a closed subset of $E_{1}$ and $W$ is an open subset of $E$ such that $W \cap X_{1} \times\{0\} \neq \emptyset$ then there exists a $C^{p}$ diffeomorphism $g$ from $E \backslash\left(X_{1} \times\{0\}\right)$ onto $E \backslash\left(\left(X_{1} \times\{0\}\right) \backslash W\right)$ such that $g$ is the identity on $(E \backslash W) \backslash\left(X_{1} \times\{0\}\right)$ and moves no point more than $\delta$.
To this end we next construct some auxiliary functions following Renz's strategy [58, pp. 54-59]. In what follows $\omega$ will denote the smooth asymmetric subadditive functional on $E_{2}$ given by Lemma 2.10 .
Lemma 2.17. There exists a continuous function $\varphi: E_{1} \rightarrow\left[0, \frac{\delta}{2}\right]$ such that:
(1) $\varphi$ is of class $C^{p}$ on $E_{1} \backslash \partial \varphi^{-1}(0)$.
(2) $W \cap\left(X_{1} \times\{0\}\right) \subset\left\{\left(x_{1}, x_{2}\right) \in E:\left\|x_{2}\right\|<\varphi\left(x_{1}\right)\right\} \subset W$.

Proof. Let $\pi_{1}: E \rightarrow E_{1}$ denote the canonical projection defined by $\pi_{1}\left(x_{1}, x_{2}\right)=x_{1}$. The set

$$
W_{1}:=\pi_{1}\left(W \cap\left(E_{1} \times\{0\}\right)\right)
$$

is open in $E_{1}$, and the function $G: E_{1} \rightarrow[0, \infty)$ defined by

$$
G\left(x_{1}\right)=\min \left\{\frac{\delta}{2}, \operatorname{dist}\left(\left(x_{1}, 0\right), E \backslash W\right)\right\}
$$

is continuous and satisfies that $G>0$ on $W_{1}$ and $G=0$ on $\pi_{1}\left(\left(E_{1} \times\{0\}\right) \backslash W\right)$. Since $E_{1}$ has $C^{p}$ smooth partitions of unity and $G$ is continuous and strictly positive on $W_{1}$, we can find a $C^{p}$ smooth function $F$ on $W_{1}$ such that

$$
0<\frac{1}{4} G\left(x_{1}\right)<F\left(x_{1}\right)<\frac{1}{2} G\left(x_{1}\right)
$$

for every $x_{1} \in W_{1}$. Now let us define $\varphi: E_{1} \rightarrow[0,1]$ by

$$
\varphi\left(x_{1}\right)= \begin{cases}F\left(x_{1}\right) & \text { if } x_{1} \in W_{1} \\ 0 & \text { if } x_{1} \in E_{1} \backslash W_{1}\end{cases}
$$

It is immediately seen that $\varphi$ is continuous, and of course $\varphi$ is of class $C^{p}$ on $E_{1} \backslash \partial \varphi^{-1}(0)=W_{1} \cup$ $\operatorname{int}\left(E_{1} \backslash W_{1}\right)$. Since $\varphi\left(x_{1}\right)=F\left(x_{1}\right)>0$ for all $x_{1} \in W_{1}$, it is obvious that

$$
W \cap\left(X_{1} \times\{0\}\right) \subset W_{1} \times\{0\} \subset\left\{\left(x_{1}, x_{2}\right) \in E:\left\|x_{2}\right\|<\varphi\left(x_{1}\right)\right\} .
$$

On the other hand, if $\left\|x_{2}\right\|<\varphi\left(x_{1}\right)$ then observe first that $x_{1} \in W_{1}$ (because if $\varphi\left(x_{1}\right)=0$ the inequality is impossible). We then must have $\left(x_{1}, x_{2}\right) \in W$, as otherwise we would get

$$
\begin{aligned}
& \operatorname{dist}\left(\left(x_{1}, 0\right), E \backslash W\right) \leq \operatorname{dist}\left(\left(x_{1}, 0\right),\left(x_{1}, x_{2}\right)\right)+\operatorname{dist}\left(\left(x_{1}, x_{2}\right), E \backslash W\right)= \\
& \left\|x_{2}\right\|+0=\left\|x_{2}\right\|<\varphi\left(x_{1}\right)<\frac{1}{2} G\left(x_{1}\right) \leq \frac{1}{2} \operatorname{dist}\left(\left(x_{1}, 0\right), E \backslash W\right),
\end{aligned}
$$

which is absurd. This shows that $\left\{\left(x_{1}, x_{2}\right) \in E:\left\|x_{2}\right\|<\varphi\left(x_{1}\right)\right\} \subset W$ and concludes the proof of the lemma.

We will also need to use a diffeomorphism $h_{2}$ of $E_{2}$ onto itself which carries the unit ball of $E_{2}$ onto the convex body $\left\{x_{2} \in E_{2}: \omega(x) \leq 1\right\}$ and such that $h_{2}(0)=0$. The existence of $h_{2}$ is ensured by the following lemma. We say that a convex body $U$ which contains 0 as an interior point is radially bounded provided that for every $x \in U$ the set $\{t x: t \in[0, \infty)\} \cap U$ is bounded.
Lemma 2.18. Let $X$ be a Banach space, and let $U_{1}, U_{2}$ be radially bounded, $C^{p}$ smooth convex such that the origin is an interior point of both $U_{1}$ and $U_{2}$. Then there exists a $C^{p}$ diffeomorphism $g: X \longrightarrow X$ such that $g\left(U_{1}\right)=U_{2}, g(0)=0$, and $g\left(\partial U_{1}\right)=\partial U_{2}$.
Proof. If $U$ and $V$ are $C^{p}$ smooth, radially bounded convex bodies such that the origin is an interior point of both $U$ and $V$, and we additionally assume that $U \subseteq V$, such a diffeomorphism can be constructed as follows: let $\theta(t)$ be a non-decreasing real function of class $C^{\infty}$ defined for $t>0$, such that $\theta(t)=0$ for $t \leq 1 / 2$ and $\theta(t)=1$ for $t \geq 1$, and define

$$
g(x)=\left(\theta\left(\mu_{U}(x)\right) \frac{\mu_{U}(x)}{\mu_{V}(x)}+1-\theta\left(\mu_{U}(x)\right)\right) x
$$

for $x \neq 0$, and $g(0)=0$. Here $\mu_{A}$ denotes the Minkowski functional of $A$. In the general case, let $U=\left\{x \in X: \mu_{U_{1}}(x)+\mu_{U_{2}}(x) \leq 1\right\}$, then $U \subseteq U_{j}$, for $j=1,2$, and there exist diffeomorphisms $g_{1}, g_{2}: X \rightarrow X$ such that $g_{j}(U)=U_{j}$ and $g_{j}(\partial U)=\partial U_{j}, j=1,2$. Then $g=g_{2} \circ g_{1}^{-1}$ does the job. See [29] for details.

The following lemma is an immediate consequence of the existence of partitions of unity in $E_{1}$.
Lemma 2.19. Suppose that $E_{1}$ is a Banach space with $C^{p}$ smooth partitions of unity, and let $X_{1}$ be a closed subset of $E_{1}$. Then there exists a continuous function $\eta: E_{1} \rightarrow[0, \infty)$ such that:
(1) $X_{1}=\eta^{-1}(0)$;
(2) $\eta$ is of class $C^{p}$ on $E_{1} \backslash X_{1}$.

We are ready to proceed with the proof of Theorem [2.15, Let $\eta$ be a function as in the statement of Lemma 2.19, and pick a diffeomorphism $h_{2}: E_{2} \rightarrow E_{2}$ such that $h_{2}(0)=0$ and

$$
h_{2}\left(\left\{x_{2} \in E_{2}:\left\|x_{2}\right\| \leq 1\right\}\right)=\left\{x_{2} \in E_{2}: \omega\left(x_{2}\right) \leq 1\right\} .
$$

Let us define

$$
A:=\varphi^{-1}((0,1]) \times E_{2}=W_{1} \times E_{2},
$$

and

$$
\Phi\left(x_{1}, x_{2}\right)=\left(x_{1}, h_{2}\left(\frac{1}{\varphi\left(x_{1}\right)} x_{2}\right)\right), \quad\left(x_{1}, x_{2}\right) \in A
$$

It is clear that $\Phi: A \rightarrow A$ is a $C^{p}$ diffeomorphism, with inverse

$$
\Phi^{-1}\left(y_{1}, y_{2}\right)=\left(y_{1}, \varphi\left(y_{1}\right) h_{2}^{-1}\left(y_{2}\right)\right),
$$

and also that

$$
\Phi\left(\left(X_{1} \times\{0\}\right) \cap A\right)=\left(X_{1} \times\{0\}\right) \cap A .
$$

Next let us define $\psi: \varphi^{-1}((0,1])=\pi_{1}(A)=W_{1} \rightarrow[0, \infty)$ by

$$
\psi\left(x_{1}\right)=\frac{\eta\left(x_{1}\right)}{\varphi\left(x_{1}\right)},
$$

and notice that $\psi$ is continuous, that

$$
\psi^{-1}(0)=\pi_{1}\left(\left(X_{1} \times\{0\}\right) \cap A\right)=X_{1} \cap W_{1},
$$

and that $\psi$ is of class $C^{p}$ outside $\psi^{-1}(0)$. By Theorem 2.8 we can find a $C^{p}$ diffeomorphism $H$ from $A \backslash\left(X_{1} \times\{0\}\right)$ onto $A$ such that $H$ is the identity outside $\left\{\left(x_{1}, x_{2}\right) \in A \backslash\left(X_{1} \times\{0\}\right): S\left(\psi\left(x_{1}\right), \omega\left(x_{2}\right)\right)<\right.$ $1\}$, where $S$ is a smooth square. Since $\Phi: A \rightarrow A$ is a $C^{p}$ diffeomorphism which takes $\left(X_{1} \times\{0\}\right) \cap A$ onto itself, we have that the composition $\Phi^{-1} \circ H \circ \Phi$ defines a $C^{p}$ diffeomorphism from $A \backslash\left(X_{1} \times\{0\}\right)$ onto $A$. Now we extend this diffeomorphism outside $A \backslash\left(X_{1} \times\{0\}\right)$ by defining $g: E \backslash\left(X_{1} \times\{0\}\right) \rightarrow E$ by

$$
g(x)= \begin{cases}\Phi^{-1} \circ H \circ \Phi(x) & \text { if } x \in A \backslash\left(X_{1} \times\{0\}\right) \\ x & \text { if } x \in E \backslash\left(A \cup\left(X_{1} \times\{0\}\right)\right)\end{cases}
$$

This mapping is clearly a bijection. Thus, in order to see that $g$ is a $C^{p}$ diffeomorphism, it is enough to see that $g$ is locally a $C^{p}$ diffeomorphism. We already know this is so for all points of $E \backslash(\partial(A \backslash$ $\left.\left.\left(X_{1} \times\{0\}\right)\right) \cup\left(X_{1} \times\{0\}\right)\right)$. Let us show that this is also true for every point $\left(x_{1}, x_{2}\right)$ of $\partial\left(A \backslash\left(X_{1} \times\right.\right.$ $\{0\})) \backslash\left(X_{1} \times\{0\}\right)$. We have $\varphi\left(x_{1}\right)=0$, and also either $\eta\left(x_{1}\right)>0$ or $\left\|x_{2}\right\|>0$. Then

$$
\lim _{A \ni\left(y_{1}, y_{2}\right) \rightarrow\left(x_{1}, x_{2}\right)} \max \left\{\frac{\eta\left(y_{1}\right)}{\varphi\left(y_{1}\right)}, \frac{\left\|y_{2}\right\|}{\varphi\left(y_{1}\right)}\right\}=\infty
$$

hence there exists a neighborhood $V$ of $\left(x_{1}, x_{2}\right)$ in $E \backslash\left(X_{1} \times\{0\}\right)$ such that

$$
\max \left\{\frac{\eta\left(y_{1}\right)}{\varphi\left(y_{1}\right)}, \frac{\left\|y_{2}\right\|}{\varphi\left(y_{1}\right)}\right\}>1 \quad \text { for all }\left(y_{1}, y_{2}\right) \in V \cap A
$$

By Lemma 2.13(2) it follows that

$$
S\left(\psi \circ \pi_{1} \circ \Phi(y), \omega \circ \pi_{2} \circ \Phi(y)\right)>1 \text { for all } y \in V \cap A \text {, }
$$

hence that $H$ is the identity on $\Phi(V \cap A)$, and consequently that $g$ is the identity on $V$, and in particular a $C^{p}$ diffeomorphism locally at $\left(x_{1}, x_{2}\right)$. Thus $g: E \backslash\left(X_{1} \times\{0\}\right) \rightarrow E$ is a $C^{p}$ diffeomorphism.
Furthermore, if $\left\|x_{2}\right\| \geq \varphi\left(x_{1}\right),\left(x_{1}, x_{2}\right) \in A \backslash\left(X_{1} \times\{0\}\right)$, then we have $\omega\left(\pi_{2} \circ \Phi\left(x_{1}, x_{2}\right)\right) \geq 1$. Hence, as above by Lemma 2.13(2), we conclude that $H\left(\Phi\left(x_{1}, x_{2}\right)\right)=\Phi\left(x_{1}, x_{2}\right)$, and it follows that $g\left(x_{1}, x_{2}\right)=$ $\left(x_{1}, x_{2}\right)$. Thus $g$ is the identity off of the set $\left\{\left(x_{1}, x_{2}\right) \in E \backslash\left(X_{1} \times\{0\}\right):\left\|x_{2}\right\|<\varphi\left(x_{1}\right)\right\}$. Since

$$
\left\{\left(x_{1}, x_{2}\right) \in E \backslash\left(X_{1} \times\{0\}\right):\left\|x_{2}\right\|<\varphi\left(x_{1}\right)\right\} \subset(E \backslash W) \backslash\left(X_{1} \times\{0\}\right),
$$

$g$ is the identity off of the set $(E \backslash W) \backslash\left(X_{1} \times\{0\}\right)$ as well.
Finally let us check that $g$ does not move any point more than $\delta$. We know that if $g$ moves a point $\left(x_{1}, x_{2}\right) \in E \backslash\left(X_{1} \times\{0\}\right)$ then $\left\|x_{2}\right\|<\varphi\left(x_{1}\right)$, and also that $g_{2}$ only moves the second coordinate $x_{2}$, that is $g\left(x_{1}, x_{2}\right)=\left(x_{1}, \pi_{2}\left(g\left(x_{1}, x_{2}\right)\right)\right)$. Hence

$$
\left\|g\left(x_{1}, x_{2}\right)-\left(x_{1}, x_{2}\right)\right\| \leq\left\|\pi_{2}\left(g\left(x_{1}, x_{2}\right)\right)-x_{2}\right\| \leq\left\|\pi_{2}\left(g\left(x_{1}, x_{2}\right)\right)\right\|+\left\|x_{2}\right\| \leq \varphi\left(x_{1}\right)+\varphi\left(x_{1}\right) \leq \delta .
$$

The proof of Theorem 2.15 is complete.
2.4. Patching local diffeomorphisms together. Throughout this section $E$ will be a Banach space. Our goal in this section is to complete the proof of Theorem 1.4.
First let us introduce the following definitions.
Definition 2.20. We will say that a subset $X$ of $E$ has the strong $C^{p}$ extraction property with respect to an open set $U$ if $X \subseteq U, X$ is relatively closed in $U$, and for every open set $V \subseteq U$ and every subset $Y \subseteq X$ relatively closed in $U$ there exists a $C^{p}$ diffeomorphism $\varphi$ from $U \backslash Y$ onto $U \backslash(Y \backslash V)$ which is the identity on $(U \backslash V) \backslash Y$. If in addition for any $\varepsilon>0$ we can require the diffeomorphism $\varphi$ not to move any point more than $\varepsilon$, we will say that $X$ has the $\varepsilon$-strong $C^{p}$ extraction property with respect to $U$.
We will also say that such a closed set $X$ has locally the strong (or the $\varepsilon$-strong) local $C^{p}$ extraction property if for every point $x \in X$ there exists an open neighborhood $U_{x}$ of $x$ such that $X \cap \overline{U_{x}}$ has the strong ( $\varepsilon$-strong, respectively) $C^{p}$ extraction property with respect to every open set $U$ with $X \cap \overline{U_{x}} \subseteq U$. (Equivalently, there exists an open neighborhood $U_{x}$ of $x$ such that $X \cap U_{x}$ has the strong ( $\varepsilon$-strong respectively) $C^{p}$ extraction property with respect to every open set $U$ for which $X \cap U_{x}$ is a relatively closed subset of $U$.)
Remark 2.21. Let $(U, W), W \subset U$, be a pair of open sets in a Banach space $E$. We say that $(U, W)$ has the strong $C^{p}$ expansion property if, for every open subsets $V$ and $U^{\prime}$ of $U, W \subset U^{\prime}$, there exists a $C^{p}$ diffeomorphism $U^{\prime} \cap V \rightarrow V$ which, by letting $\varphi(x)=x$ for $x \in U^{\prime} \backslash V$, extends to $C^{p}$ diffeomorphism $\varphi: U^{\prime} \rightarrow U^{\prime} \cup V$.
In particular, letting $U^{\prime}=W$, there exists a $C^{p}$ diffeomorphism $W \cap V \rightarrow V$ which extends to a $C^{p}$ diffeomorphism of $W$ onto $W \cup V$ via the identity off $W \cap V$. Hence, $W$ is smoothly expanded to $W \cup V$; this justifies the term of $C^{p}$ expansion. Should this expansion be valid for all open sets $U^{\prime}$, $W \subset U^{\prime}$, then we would have the strong $C^{p}$ expansion property.
Notice that a relatively closed subset $X$ has the strong $C^{p}$ extraction property with respect to $U$ if and only if $(U, W)=(U, U \backslash X)$ has the strong $C^{p}$ expansion property.
We say that an open subset $W$ of $E$ has locally the strong $C^{p}$ expansion property if every $x \in E \backslash W$ has an open neighborhood $U_{x}$ such that $\left(U, U_{x} \cap W\right)$ has the strong expansion property for every open set $U \supset U_{x} \cap W$.
Notice that a closed set $X$ has locally the strong $C^{p}$ extraction property if and only if $W=E \backslash X$ has locally the strong $C^{p}$ expansion property.

Some basic properties that can be derived from Definition 2.20 are listed in the following lemma.

Lemma 2.22. Let us suppose that $X \subset E$ has the $\varepsilon$-strong $C^{p}$ extraction property with respect to an open set $U$ of $E$. Then:
(1) For every closed set $Y \subseteq X, Y$ has the $\varepsilon$-strong $C^{p}$ extraction property with respect to $U$;
(2) For every open subset $U^{\prime} \subseteq U, X \cap U^{\prime}$ has the $\varepsilon$-strong $C^{p}$ extraction property with respect to $U^{\prime}$.
(3) If $h$ is a $C^{p}$ diffeomorphism defined on $U$ and such that $h(U)$ is open, then $h(X)$ has the strong $C^{p}$ extraction property with respect to $h(U)$.
Proof.
(1) This follows directly from the definition.
(2) Take an open subset $V^{\prime} \subseteq U^{\prime}$, a subset $Y \subseteq X \cap U^{\prime}$ relatively closed in $U^{\prime}$. Since $X$ has the strong $C^{p}$ extraction property with respect to $U$ there exists a $C^{p}$ diffeomorphism $\varphi$ from $U \backslash \bar{Y}$ onto $U \backslash\left(\bar{Y} \backslash V^{\prime}\right)$ which is the identity on $\left(U \backslash V^{\prime}\right) \backslash \bar{Y}$. When restricting $\varphi$ to $U^{\prime} \backslash Y$ we actually get a $C^{p}$ diffeomorphism from $U^{\prime} \backslash Y$ onto $U^{\prime} \backslash\left(Y \backslash V^{\prime}\right)$ which is the identity on $\left(U^{\prime} \backslash V^{\prime}\right) \backslash Y$.
(3) Take an open subset $V$ of $h(U)$, a subset $Y \subseteq h(X)$ relatively closed in $h(U)$. Since $X$ has the strong $C^{p}$ extraction property with respect to $U$ and $h^{-1}(Y)$ is relatively closed in $U$, there exists a $C^{p}$ diffeomorphism $\varphi$ from $U \backslash h^{-1}(Y)$ onto $U \backslash\left(h^{-1}(Y) \backslash h^{-1}(V)\right)$ which is the identity on $\left(U \backslash h^{-1}(V)\right) \backslash h^{-1}(Y)$. Then the mapping

$$
g:=h \circ \varphi \circ h^{-1}: h(U) \backslash Y \longrightarrow h(U) \backslash(Y \backslash V)
$$

is a surjective $C^{p}$ diffeomorphism which restricts to the identity on $(h(U) \backslash V) \backslash Y$.
Remark 2.23. In Lemma 2.22 (3) we do not have in general the $\varepsilon$-strong $C^{p}$ extraction property of $h(X)$ with respect to $h(U)$, but we still have the following: suppose $h$ does not move any points more than some $\varepsilon>0$. For every $\eta>0$ in the proof of Lemma 2.22 (3) we can assume that $\varphi$ does not move any point more than $\eta$. Hence $g \circ h=\left(h \circ \varphi \circ h^{-1}\right) \circ h=h \circ \varphi$ does not move any point more than $\varepsilon+\eta$.

Let us state the main result of this section, which is crucial in the proof of Theorem 1.4 provided underneath.

Theorem 2.24. Let $E$ be a Banach space and $X$ be a closed subset of $E$ which has locally the $\varepsilon$-strong $C^{p}$ extraction property. Let $U$ be an open subset of $E$ and $\mathcal{G}=\left\{G_{r}\right\}_{r \in \Omega}$ be an open cover of $E$. Then there exists a $C^{p}$ diffeomorphism $g$ from $E \backslash X$ onto $E \backslash(X \backslash U)$ which is the identity on $(E \backslash U) \backslash X$ and is limited by $\mathcal{G}$.
Proof of Theorem 1.4. Let us show that $X$ from the statement of Theorem 1.4 has locally the strong $C^{p}$ extraction property.
To this end, fix $x \in X$ and choose a neighborhood $U_{x}$ such that $X \cap U_{x} \subset G\left(f_{x}\right)$; we can assume that $U_{x}=\overline{U_{x}}$. Further, we can assume that $f_{x}$ is defined and continuous on the whole $E_{(1, x)}$. We will show that $X^{\prime}:=X \cap \overline{U_{x}}$ has the $\varepsilon$-strong $C^{p}$ extraction property with respect to every open set $U$ with $X^{\prime} \subseteq U$. Notice that $X^{\prime}=G\left(f_{x} \mid X_{1}^{\prime}\right)$ for a certain closed $X_{1}^{\prime} \subset E_{(1, x)}$. Furthermore, if $Y^{\prime} \subset X^{\prime}$ is relatively closed in $U$, then $Y^{\prime}$ is closed in $X^{\prime}$. Hence, $Y^{\prime}=G\left(f_{x} \mid Y_{1}^{\prime}\right)$ for a certain closed $Y_{1}^{\prime} \subset X_{1}^{\prime} \subset E_{(1, x)}$. Let $V$ be an open subset of $U$. Take now $\varepsilon>0$ and apply Theorem 2.15 to $E_{1}:=E_{(1, x)}, E_{2}:=E_{(2, x)}, f:=f_{x}, X:=Y^{\prime}$, and $V$ (in place of $U$ ) to obtain a $C^{p}$ diffeomorphism $g$ from $E \backslash Y^{\prime}$ onto $E \backslash\left(Y^{\prime} \backslash V\right)$ such that $g$ is the identity on $(E \backslash V) \backslash Y^{\prime}$ and moves no point more than $\varepsilon$. Then $\varphi=g \mid U$ is as required in the definition of the $\varepsilon$-strong $C^{p}$ extraction of $X^{\prime}$ with respect to $U$.
Now, an application of Theorem 2.24 concludes our proof.

The remaining part of this section is devoted to proving Theorem 2.24. Firstly, the fact that we are working with a set $X$ that has locally the strong $C^{p}$ extraction property, and the requirement that our final $C^{p}$ diffeomorphism must be limited by a given open cover $\mathcal{G}$ forces us to employ good refinements of covers of the Banach space $E$. In the separable case star-finite refinements provide an adequate tool to face the problem (see West [67). Recall that a cover is said to be star-finite provided that each element of the cover intersects at most finitely many others. However, in the nonseparable case, getting a star-finite refinement of an open cover, in general, is not possible. We will use sigma-discrete refinements as shown in the following.
Lemma 2.25. Let $E$ be a Banach space and $X$ be a closed subset of $E$ which has locally the $\varepsilon$-strong $C^{p}$ extraction property. Let $\mathcal{G}=\left\{G_{r}\right\}_{r \in \Omega}$ be an open cover of $E$, where the cardinality of the indexing set $\Omega$ is the density of $E$. Then there exist collections $\left\{X_{i}\right\}_{i \geq 1},\left\{W_{i}\right\}_{i \geq 1},\left\{V_{i}\right\}_{i \geq 1}$, such that:
(1) $X_{i} \subseteq W_{i} \subseteq \overline{W_{i}} \subseteq V_{i}$ for all $i \in \mathbb{N}$;
(2) $\left\{V_{i}\right\}_{i \geq 1}$ and $\left\{W_{i}\right\}_{i \geq 1}$ are star-finite open covers of $E$;
(3) $\left\{X_{i}\right\}_{i \geq 1}$ is a cover of $X$ by closed subsets of $X$;
(4) Each $\bar{W}_{i}$ and $V_{i}$ admits an open discrete cover $\left\{W_{i, r}\right\}_{r \in \Omega}$ and $\left\{V_{i, r}\right\}_{r \in \Omega}$, repectively; more precisely,

$$
\begin{gathered}
W_{i}=\bigcup_{r \in \Omega} W_{i, r} \quad \text { and } \quad V_{i}=\bigcup_{r \in \Omega} V_{i, r}, \\
\overline{W_{i, r}} \subseteq V_{i, r} \text { for every } r \in \Omega,
\end{gathered}
$$

and

$$
\operatorname{dist}\left(V_{i, r}, V_{i, r^{\prime}}\right) \geq \frac{1}{2^{i+1}} \text { for every } r, r^{\prime} \in \Omega, r \neq r^{\prime}
$$

(5) $\left\{W_{i, r}\right\}_{i \geq 1, r \in \Omega}$ and $\left\{V_{i, r}\right\}_{i \geq 1, r \in \Omega}$ are open refinements of $\mathcal{G}$;
(6) Each $X_{i}$ can be written as $X_{i}=\bigcup_{r \in \Omega} X_{i, r}$, where $X_{i, r}$ is a closed subset of $X$ satisfying the following requirements

$$
X_{i, r} \subseteq W_{i, r} \subseteq \overline{W_{i, r}} \subseteq V_{i, r}
$$

and

$$
X_{i, r} \text { has the } \varepsilon-\text { strong } C^{p} \text { extraction property with respect to } V_{i, r} \text {. }
$$

Proof. For each $x \in E$, let $U_{x}$ be an open neighborhood of $x$ such that $\overline{U_{x}} \subseteq G$ for some $G \in \mathcal{G}$ and also satisfying that
(1) if $x \in X$ then $X \cap \overline{U_{x}}$ has the $\varepsilon$-strong $C^{p}$ extraction property with respect to every open set $U$ with $X \cap \overline{U_{x}} \subseteq U$, and
(2) if $x \notin X$ then $X \cap \overline{U_{x}}=\emptyset$.

Since the cardinality of $\Omega$ is the density of $E$, we can extract a subcover $\mathcal{U}=\left\{U_{r}: r \in \Omega\right\}$ from $\left\{U_{x}: x \in E\right\}$. Now we use a result of Rudin [61] (see also [38, p. 390]) to obtain two open refinements $\left\{A_{j, r}\right\}_{j \geq 1, r \in \Omega}$ and $\left\{B_{j, r}\right\}_{j \geq 1, r \in \Omega}$ of $\mathcal{U}$ such that
(1) $A_{j, r} \subseteq B_{j, r} \subseteq U_{r}$ for all $j \in \mathbb{N}$ and $r \in \Omega$;
(2) $\operatorname{dist}\left(A_{j, r}, E \backslash B_{j, r}\right) \geq \frac{1}{2^{j}}$ for all $j \in \mathbb{N}$ and $r \in \Omega$;
(3) $\operatorname{dist}\left(B_{j, r}, B_{j, r^{\prime}}\right) \geq \frac{1}{2^{j+1}}$ for all $j \in \mathbb{N}$ and $r, r^{\prime} \in \Omega, r \neq r^{\prime}$;
(4) Letting $A_{j}=\bigcup_{r \in \Omega} A_{j, r}$ and $B_{j}=\bigcup_{j \in \Omega} B_{j, r}$ each collection $\left\{A_{j}\right\}_{j \geq 1},\left\{B_{j}\right\}_{j \geq 1}$ forms a locally finite open cover of $E$.
Observe that $\overline{A_{j}} \subseteq B_{j}$ for every $j \in \mathbb{N}$.
For every $j$, there exists a sequence of open sets $B_{j}^{n}, n \geq j$, so that

$$
\overline{A_{j}} \subset B_{j}^{j} \subset \overline{B_{j}^{j}} \subset B_{j}^{j+1} \subset \overline{B_{j}^{j+1}} \subset \cdots \subset B_{j}^{n} \subset \overline{B_{j}^{n}} \subset B_{j}^{n+1} \subset \cdots \subset B_{j} .
$$

For each $j$, write $\mathcal{B}_{j}=\left\{B_{j}^{n}: n \geq j\right\}$. Clearly, $\mathcal{B}=\bigcup_{j=1}^{\infty} \mathcal{B}_{j}$ is an open cover of $E$; likewise, the family $\{\bar{B} \cap X: B \in \mathcal{B}\}$ is a closed cover of $X$.
Defining for each $n \in \mathbb{N}: Y_{n}:=\bigcup_{j=1}^{n} B_{j}^{n}, H_{n}:=Y_{n} \backslash \overline{Y_{n-3}}$ and $K_{n}:=\overline{Y_{n}} \backslash Y_{n-1}$ (let $Y_{-2}=Y_{-1}=Y_{0}=$ $\emptyset$ ), we have the following properties:

- $E=\bigcup_{n=1}^{\infty} Y_{n}$;
- $\overline{Y_{n}} \subseteq Y_{n+1}$ for all $n \in \mathbb{N}$;
- $K_{n} \subseteq H_{n+1}$ for all $n \in \mathbb{N}$;
- $E=\bigcup_{n=1}^{\infty} K_{n}$;
- $H_{m} \cap H_{n}=\emptyset$ for all $m, n$ with $|m-n| \geq 3$.

Hence, the collection

$$
\bigcup_{n=1}^{\infty}\left\{K_{n} \cap \overline{B_{j}^{n}}: j=1, \ldots, n\right\}
$$

is a closed cover of $E$ and therefore

$$
\bigcup_{n=1}^{\infty}\left\{H_{n+1} \cap B_{j}^{n+1}: j=1, \ldots, n\right\}
$$

is an open cover of $E$. Both covers are countable and star-finite, and they are refinements of $\left\{B_{j}\right\}_{j \geq 1}$. We call the first one $\left\{T_{i}\right\}_{i \geq 1}$ and the second one $\left\{V_{i}\right\}_{i \geq 1}$, that is, for every $i$ there corresponds a unique pair $(j, n), n \geq j$, with $T_{i}=K_{n} \cap \overline{B_{j}^{n}}$ and $V_{i}=H_{n+1} \cap B_{j}^{n+1}$. Consequently, we have $T_{i} \subseteq V_{i}$ for every $i \in \mathbb{N}$.
Now for each $i \in \mathbb{N}$ we take $j=j(i) \in \mathbb{N}$ such that $T_{i} \subseteq V_{i} \subseteq B_{j}$. Let us assume without loss of generality that $j(i) \leq i$. We can write

$$
T_{i}=\bigcup_{r \in \Omega} T_{i} \cap B_{j, r} \text { and } V_{i}=\bigcup_{r \in \Omega} V_{i} \cap B_{j, r}
$$

and we define $T_{i, r}=T_{i} \cap B_{j, r}$ and $V_{i, r}=V_{i} \cap B_{j, r}$ for every $i \in \mathbb{N}$ and $r \in \Omega$. Clearly we have that $T_{i, r} \subseteq V_{i, r}$ for all $i \in \mathbb{N}$ and $r \in \Omega$. Also dist $\left(V_{i, r}, V_{i, r^{\prime}}\right) \geq \frac{1}{2^{j+1}} \geq \frac{1}{2^{2+1}}$ for all $i \in \mathbb{N}$ and $r, r^{\prime} \in \Omega, r \neq r^{\prime}$. Finally let us define $X_{i, r}=X \cap T_{i, r}$. Bearing in mind that $T_{i, r} \subset V_{i, r} \subset B_{j, r} \subset U_{x}$ for some $x \in X$ and that $T_{i, r}$ is closed, we obtain $X_{i, r}=X \cap T_{i, r} \subset X \cap \overline{\bar{U}_{x}}$. Since $X \cap \overline{U_{x}}$ has the $\varepsilon$-strong $C^{p}$ extraction property with respect to every open set $U$ with $X \cap \overline{U_{x}} \subset U$, applying Lemma 2.22(1), we get that $X_{i, r}$ has the $\varepsilon$-strong $C^{p}$ extraction property with respect to every such an open set $U$. Finally, applying Lemma [2.22(2), $X_{i, r}$ has the strong $C^{p}$ extraction property with respect to every open set $U^{\prime}$ with $X_{i, r} \subset U^{\prime}$. In particular, $X_{i, r}$ has the $\varepsilon$-strong $C^{p}$ extraction property with respect to $V_{i, r}$.

Let us consider now for each $i \in \mathbb{N}$ and $r \in \Omega$ an open set $W_{i, r}$ with $T_{i, r} \subseteq W_{i, r} \subseteq \overline{W_{i, r}} \subseteq V_{i, r}$ and call $W_{i}=\bigcup_{r \in \Omega} W_{i, r}$. We still have that $\left\{W_{i, r}\right\}_{i \geq 1, r \in \Omega}$ is a refinement of $\mathcal{G}$, and that $\left\{W_{i}\right\}_{i \geq 1}$ is a star-finite open cover of $E$.
Then the collections $\left\{X_{i, r}\right\}_{i \geq 1, r \in \Omega},\left\{W_{i, r}\right\}_{i \geq 1, r \in \Omega},\left\{V_{i, r}\right\}_{i \geq 1, r \in \Omega}$ have the required properties. Note that each $X_{i, r}$ has the strong $C^{p}$ extraction property with respect to every open set containing it.

Lemma 2.26. Let $X_{i}, V_{i}$, and $V_{i, r}$ be as in Lemma 2.25. Then, for every $i, j \in \mathbb{N}, X_{i} \cap V_{j}$, has the strong $C^{p}$ extraction property with respect to $V_{j}$. Moreover, if $\varphi_{i, j}$ is a $C^{p}$ diffeomorphism satisfying the definition of the $\varepsilon$-strong extraction property for these sets, then $\varphi\left(V_{i, r}\right) \subset V_{i, r}$ and $\varphi\left(V_{j, r}\right) \subset V_{j, r}$ for every $r \in \Omega$.

Proof. Take $V \subseteq V_{j}$ an open set. Let $r \in \Omega$ and consider the open set $V_{j, r}$. For each $s \in \Omega$ the set $X_{i, s} \cap V_{j, r}$ has the $\varepsilon$-strong $C^{p}$ extraction property with respect to the open set $V_{i, s} \cap V_{j, r}$. We have

$$
X_{i, s} \cap V_{j, r} \subseteq W_{i, s} \cap V_{j, r} \subseteq \overline{W_{i, s}} \cap V_{j, r} \subseteq V_{i, s} \cap V_{j, r}
$$

There exists a $C^{p}$ diffeomorphism $h_{r, s}:\left(V_{i, s} \cap V_{j, r}\right) \backslash X_{i, s} \rightarrow V_{i, s} \cap V_{j, r} \backslash\left(X_{i, s} \backslash\left(W_{i, s} \cap V\right)\right)$ which is the identity on $\left(\left(V_{i, s} \cap V_{j, r} \backslash\left(W_{i, s} \cap V\right)\right) \backslash X_{i, s}\right.$. Outside $V_{i, s} \cap V_{j, r}$ we define $h_{r, s}$ to be the identity. Since $\overline{W_{i, s}} \subseteq V_{i, s}$ we have a well-defined $C^{p}$ diffeomorphism

$$
h_{r, s}: V_{j, r} \backslash X_{i, s} \rightarrow V_{j, r} \backslash\left(X_{i, s} \backslash\left(W_{i, s} \cap V\right)\right)
$$

which is the identity on $\left(V_{j, r} \backslash\left(W_{i, s} \cap V\right)\right) \backslash X_{i, s}$. In particular $h_{r, s}$ is the identity on $V_{i, s^{\prime}} \cap V_{j, r}$ for every $s^{\prime} \in \Omega, s \neq s^{\prime}$.
Having defined $h_{r, s}$ for each $s \in \Omega$ in the way described above, we finally define

$$
h_{r}=\bigcirc_{s \in \Omega} h_{r, s}
$$

as an infinite composition of $h_{r, s}, s \in \Omega$. It is easy to see that $h_{r}$ is well defined and provides a $C^{p}$ diffeomorphism of $V_{j, r} \backslash X_{i}$ onto $V_{j, r} \backslash\left(X_{i} \backslash V\right)$ which is the identity on $\left(V_{j, r} \backslash V\right) \backslash X_{i}$.
By the discreteness of the family $\left\{V_{j, r}\right\}_{r \in \Omega}$, the formula $h(x)=h_{r}(x), x \in V_{j, r} \backslash X_{i}$, defines a $C^{p}$ diffeomorphism of $\bigcup_{r \in \Omega} V_{j, r} \backslash X_{i}=V_{j} \backslash X_{i}$ onto $\bigcup_{r \in \Omega}\left(V_{j, r} \backslash\left(X_{i} \backslash V\right)\right)=V_{j} \backslash\left(X_{i} \backslash V\right)$ which is the identity on $\bigcup_{r \in \Omega}\left(\left(V_{j, r} \backslash V\right) \backslash X_{i}\right)=\left(V_{j} \backslash V\right) \backslash X_{i}$.
To end the proof observe that each $h_{r, s}$ is the identity outside $V_{i, s} \cap V_{j, r}$, so $h$ sends $V_{i, r}$ into $V_{i, r}$ and $V_{j, r}$ into $V_{j, r}$ for every $r \in \Omega$.
Notice that in the previous two Lemmas 2.26 and 2.25 one can replace the $\varepsilon$-strong $C^{p}$ extraction property with just the strong $C^{p}$ extraction property.
The last tool that we need to introduce before going into the proof of Theorem 2.24 is the next lemma (see Statement A in West's paper [67, pp. 289]).

Lemma 2.27. Let $V_{0}, \ldots, V_{n}$ be open sets of $E$, and $X_{0}, \ldots, X_{n}$ be subsets of $V_{0}, \ldots, V_{n}$. Take also an open set $U$. Suppose each $X_{i}$ is relatively closed in $V=\bigcup_{i=0}^{n} V_{i}$ and has the $\varepsilon$-strong $C^{p}$ extraction property with respect to $V$. Then for every $\varepsilon>0$ there exists a $C^{p}$ diffeomorphism of $V \backslash X_{0}$ onto $V \backslash\left(X_{0} \backslash U\right)$ which is the identity outside $V_{0} \cap U$, carries $X_{i} \backslash X_{0}$ into $V_{i}$ for each $i=1, \ldots, n$, and moves no point more than $\varepsilon$.
Proof. We will divide the set $X_{0}$ that we want to extract as follows. For each $j=0, \ldots, n$, let

$$
Q_{j}:=\left\{Z: Z=\bigcap_{i=0}^{n-j} X_{p(i)} \backslash \bigcup_{i=n-j+1}^{n} X_{p(i)} \text { for some permutation } p \text { of }\{0, \ldots, n\} \text { carrying } 0 \text { to } 0\right\} ;
$$

these are families of subsets of $X_{0}$. Let $Q=\bigcup_{j=0}^{n} Q_{j}$. The family $Q$ is a pairwise disjoint cover of $X_{0}$ with cardinality $\leq 2^{n}$. Order $Q$ in such a manner that if $j<k$ then all elements of $Q_{j}$ precede those of $Q_{k}$. We then will list the elements of $Q$ as $Z_{1}, Z_{2}, \ldots, Z_{2^{n}}$ (bearing in mind that some $Z_{m}$ 's may repeat in that listing); that is, $Q=\left\{Z_{m}\right\}_{m=1}^{2^{n}}$ and $\bigcup_{i=1}^{2^{n}} Z_{i}=X_{0}$. Note that $Q_{0}=\left\{Z_{1}\right\}$ and $Q_{n}=\left\{Z_{2^{n}}\right\}$, where

$$
Z_{1}=\bigcap_{i=0}^{n} X_{i} \quad \text { and } \quad Z_{2^{n}}=X_{0} \backslash \bigcup_{i=1}^{n} X_{i} .
$$

Likewise, for each $m=1, \ldots, 2^{n}$, if $Z_{m} \in Q_{j}$ and $p$ is a permutation for which $Z_{m}=\bigcap_{i=0}^{n-j} X_{p(i)} \backslash$ $\bigcup_{i=n-j+1}^{n} X_{p(i)}$, we define

$$
N_{m}=\bigcap_{i=0}^{n-j} V_{p(i)} \subset V_{0} .
$$

The family $\left\{N_{m}\right\}_{m=1}^{2^{n}}$ is an open cover of $V_{0}$; we also have $N_{1}=\bigcap_{i=0}^{n} V_{i}$ and $N_{2^{n}}=V_{0}$. Denote by $Q_{k}^{*}$ the union of all the elements of $Q_{k}, k=0,1, \ldots, n$; notice that $Q_{0}^{*}=Z_{1}$ and $Q_{n}^{*}=Z_{2^{n}}$. For each $j>0$, the elements of $Q_{j}$ form a pairwise disjoint family of relatively closed subsets of the open set $V \backslash \bigcup_{k=0}^{j-1} Q_{k}^{*}$. Also each $Z_{m}$ in $Q_{j}$ lies in $N_{m}$. Therefore, for each $j>0$, there exists a collection of pairwise disjoint open sets $M_{m}$ in $V \backslash \bigcup_{k=0}^{j-1} Q_{k}^{*}$, one for each $Z_{m}$ in $Q_{j}$ (that is, if $Z_{m}=Z_{m^{\prime}} \in Q_{j}$, $m \neq m^{\prime}$, then $M_{m}=M_{m^{\prime}}$, such that

$$
Z_{m} \subseteq M_{m} \subseteq N_{m} \backslash \bigcup_{i=n-j+1}^{n} X_{p(i)}
$$

where $j$ is such that $Z_{m}$ is in $Q_{j}$ and $p$ is a permutation defining $Z_{m}$ as above.
The set $Z_{1}=\bigcap_{i=0}^{n} X_{i} \subseteq N_{1}=\bigcap_{i=0}^{n} V_{i}$ is relatively closed in $V$ and has the $\varepsilon$-strong $C^{p}$ extraction property with respect to $V$, so there is a $C^{p}$ diffeomorphism

$$
h_{1}: V \backslash Z_{1} \rightarrow V \backslash\left(Z_{1} \backslash U\right)
$$

which is the identity outside $\left(N_{1} \cap U\right) \cup Z_{1}$ and moves no point more than $\frac{\varepsilon}{2^{n}}$; in particular, $h_{1}\left(Z_{2}\right) \backslash U=$ $Z_{2} \backslash U$ should $Z_{2} \neq Z_{1}$.
We will apply induction to prove that for $1<m \leq 2^{n}$ there exists a $C^{p}$ diffeomorphism

$$
h_{m}:\left(V \backslash\left(\bigcup_{i=1}^{m-1} Z_{i} \backslash U\right)\right) \backslash h_{m-1} \circ \cdots \circ h_{1}\left(Z_{m}\right) \rightarrow V \backslash\left(\bigcup_{i=1}^{m} Z_{i} \backslash U\right)
$$

which is the identity outside $\left(h_{m-1} \circ \cdots \circ h_{1}\left(M_{m}\right) \cap N_{m} \cap U\right) \cup h_{m-1} \circ \cdots \circ h_{1}\left(Z_{m}\right)$, satisfies
$h_{m} \circ h_{m-1} \circ \cdots \circ h_{1}\left(V \backslash \bigcup_{i=1}^{m} Z_{i}\right)=V \backslash\left(\bigcup_{i=1}^{m} Z_{i} \backslash U\right)$ and $h_{m} \circ h_{m-1} \circ \cdots \circ h_{1}\left(Z_{m+1}\right) \backslash U=Z_{m+1} \backslash U$,
and such that $h_{m} \circ \cdots \circ h_{1}$ moves no point more than $\frac{m \varepsilon}{2^{n}}$.
Suppose this is true for every $1 \leq k \leq m-1$, and let us check so it is for $m$. Assume $Z_{m} \in Q_{j}$; additionally, we can assume that $Z_{m-1} \neq Z_{m}$, otherwise, the identity in place of $h_{m}$ will do. By the definition of the $\varepsilon$-strong $C^{p}$ extractibility and Lemma [2.22(2), $Z_{m}$ has the $\varepsilon$-strong $C^{p}$ extraction property with respect to $V \backslash \bigcup_{k=0}^{j-1} Q_{k}^{*}$. Furthermore, one more application of Lemma [2.22(2) yields that $Z_{m} \subset V \backslash \bigcup_{i=1}^{m-1} Z_{i}$ is relatively closed and has the $\varepsilon$-strong $C^{p}$ extraction property with respect to $V \backslash \bigcup_{i=1}^{m-1} Z_{i}$. By Lemma 2.22 (3), $h_{m-1} \circ \cdots \circ h_{1}\left(Z_{m}\right) \subset h_{m-1} \circ \cdots \circ h_{1}\left(V \backslash \bigcup_{i=1}^{m-1} Z_{i}\right)$ is relatively closed and has the strong $C^{p}$ extraction property with respect to the open set

$$
h_{m-1} \circ \cdots \circ h_{1}\left(V \backslash \bigcup_{i=1}^{m-1} Z_{i}\right)=V \backslash\left(\bigcup_{i=1}^{m-1} Z_{i} \backslash U\right)
$$

Considering the open set $h_{m-1} \circ \cdots \circ h_{1}\left(M_{m}\right) \cap N_{m} \cap U$, there exists a $C^{p}$ diffeomorphism $h_{m}$ from

$$
\left(V \backslash\left(\bigcup_{i=1}^{m-1} Z_{i} \backslash U\right)\right) \backslash h_{m-1} \circ \cdots \circ h_{1}\left(Z_{m}\right)
$$

onto

$$
\begin{equation*}
\left(V \backslash\left(\bigcup_{i=1}^{m-1} Z_{i} \backslash U\right)\right) \backslash\left(h_{m-1} \circ \cdots \circ h_{1}\left(Z_{m}\right) \backslash\left(h_{m-1} \circ \cdots \circ h_{1}\left(M_{m}\right) \cap N_{m} \cap U\right)\right. \tag{2.9}
\end{equation*}
$$

which is the identity outside $\left(h_{m-1} \circ \cdots \circ h_{1}\left(M_{m}\right) \cap N_{m} \cap U\right) \cup h_{m-1} \circ \cdots \circ h_{1}\left(Z_{m}\right)$. Furthermore by Remark 2.23 and the fact that $h_{m-1} \circ \cdots \circ h_{1}$ moves no point more than $\frac{(m-1) \varepsilon}{2^{n}}$ (induction hypothesis), we have that $h_{m} \circ\left(h_{m-1} \circ \cdots \circ h_{1}\right)$ moves no point more than $\frac{\varepsilon}{2^{n}}+\frac{(m-1) \varepsilon}{2^{n}}=\frac{m \varepsilon}{2^{n}}$.
Finally note that the expression (2.9) is equal to $V \backslash\left(\bigcup_{i=1}^{m} Z_{i}^{2} \backslash U\right)$ because $h_{m-1} \circ \cdots \circ h_{1}\left(Z_{m}\right) \subseteq$ $h_{m-1} \circ \cdots \circ h_{1}\left(M_{m}\right) \cap N_{m}$ and the fact that $h_{m-1} \circ \cdots \circ h_{1}\left(Z_{m}\right) \backslash U=Z_{m} \backslash U$.

To conclude define

$$
h=h_{2^{n}} \circ \cdots \circ h_{1}: V \backslash X_{0} \rightarrow V \backslash\left(X_{0} \backslash U\right) .
$$

This is a $C^{p}$ diffeomorphism of $V \backslash X_{0}$ onto $V \backslash\left(X_{0} \backslash U\right)$ which is the identity outside $\left(V_{0} \cap U\right) \cup X_{0}$ and moves no point more than $\varepsilon$.
To end the proof we will show that $h$ carries $X_{i} \backslash X_{0}$ into $V_{i}$ for each $i=1, \ldots, n$. Take $x \in X_{i} \backslash X_{0}$ for some $i=1, \ldots, n$ and let us see that $h(x) \in V_{i}$. If $h_{1}(x) \neq x$ then $h_{1}(x) \in N_{1}$, so because $N_{1} \subseteq V_{i}$ we have that $h_{1}(x) \in V_{i}$. Suppose now that $h_{m-1} \circ \cdots \circ h_{1}(x) \in V_{i}$. If $h_{m} \circ \cdots \circ h_{1}(x) \neq$ $h_{m-1} \circ \cdots \circ h_{1}(x)$ then we have that $x \in M_{m}$ and $h_{m} \circ \cdots \circ h_{1}(x) \in N_{m}$. We have that $x \in X_{i} \cap M_{m}$ and $M_{m} \subseteq N_{m} \backslash \bigcup_{i=n-j+1}^{n} X_{p(i)}$ where $j$ is such that $Z_{m}$ is in $Q_{j}$ and $p$ is a permutation that defines $Z_{m}$. Obviously we must have that $i=p\left(i_{0}\right)$ where $i_{0}=\{1, \ldots, n-j\}$. Also $N_{m}=\bigcap_{k=0}^{n-j} V_{p(k)} \subseteq V_{i}$, hence $h_{m} \circ \cdots \circ h_{1}(x)$ must also lie in $V_{i}$. Applying induction we have that $h(x) \in V_{i}$.

Remark 2.28. Notice that if we only have the strong $C^{p}$ extraction property of the sets $X_{i}$ with respect to $V$, we get the same result except that for any $\varepsilon>0$ we can not assure that the final extracting diffeomorphism moves all points less than $\varepsilon$. In such a case we will say that we have a weak version of Lemma 2.27.
However, suppose we have that $X_{0} \subseteq V_{0}, X_{1} \subseteq V_{1}, \ldots, X_{n} \subseteq V_{n}$ are of the form $g\left(X_{i}\right) \subseteq g\left(V_{i}\right)$, $i=0,1, \ldots, n$, where $g$ is a $C^{p}$ diffeomorphism that moves points less than some $\delta_{1}>0$. In such a case, using Remark 2.23, for any $\delta_{2}>0$ we can make the final extracting diffeomorphism $h$ of the proof of Lemma 2.27 satisfy that $h \circ g$ moves no point more than $\delta_{1}+\delta_{2}$. In particular if $g(x)=x$ then we have that $\|h(g(x))-x\| \leq \delta_{2}$.

Remark 2.29. Assume, additionally, that the set $V$ of Lemma 2.27 is of the form $V=\bigcup_{r \in \Omega} V_{r}$, where $\Omega$ is a set of indexes, each $V_{r}$ is an open set, and $V_{r} \cap V_{r^{\prime}}=\emptyset$ for every $r, r^{\prime} \in \Omega, r \neq r^{\prime}$. Then we can also require that the extracting $C^{p}$ diffeomorphism of $V \backslash X_{0}$ onto $V \backslash\left(X_{0} \backslash U\right)$ sends each set $V_{r} \backslash X_{0}$ into $V_{r}$ for every $r \in \Omega$. To prove this, fix $r \in \Omega$ and replace the sequences $V_{0}, V_{1}, \ldots, V_{n}$ and $X_{0}, X_{1}, \ldots X_{n}$ with $V_{0} \cap V_{r}, V_{1} \cap V_{r} \ldots, V_{n} \cap V_{r}$ and $X_{0} \cap V_{r}, X_{1} \cap V_{r}, \ldots, V_{n} \cap V_{r}$, respectively. Further, observe that $\bigcup_{i=0}^{n} V_{i} \cap V_{r}=V_{r}$ and that each $X_{i} \cap V_{r}$ has the strong $C^{p}$ extraction property with respect to $V_{r}$ by Lemma 2.22 (2). According to the assertion of Lemma 2.27, we conclude that there exists a $C^{p}$ diffeomorphism $h_{r}: V_{r} \backslash X_{0} \rightarrow V_{r} \backslash\left(X_{0} \backslash U\right)$ satisfying the required conditions. Finally, it is enough to set $h: V \backslash X_{0} \rightarrow V \backslash\left(X_{0} \backslash U\right)$ by letting $h(x)=h_{r}(x)$ for $x \in V_{r} \backslash X_{0}$.

The rest of the proof of Theorem [2.24 goes as in [67, Theorem 1], with some modifications due to the facts that we work here with an open set $U$ not necessarily containing $X$, and that $E$ is not necessarily separable.

Proof of Theorem 2.24. Apply Lemma 2.25 to the given cover $\mathcal{G}$ to find collections $\left\{X_{i}\right\}_{i>1},\left\{W_{i}\right\}_{i \geq 1}$ and $\left\{V_{i}\right\}_{i \geq 1}$ of subsets of $E$ satisfying conditions (1)-(6) of that lemma. By Lemma [2.26, for all $i, j \in \mathbb{N}, X_{i} \cap V_{j}$ has the $\varepsilon$-strong $C^{p}$ extraction property with respect to $V_{j}$. Moreover, if $\varphi_{i, j}$ is a $C^{p}$ diffeomorphism with this property, then $\varphi\left(V_{i, r}\right) \subset V_{i, r}$ and $\varphi\left(V_{j, r}\right) \subset V_{j, r}$ for every $r \in \Omega$.
Let us now define the required $C^{p}$ diffeomorphism $g: E \backslash X \rightarrow E \backslash(X \backslash U)$ which is the identity on $(E \backslash U) \backslash X$ and is limited by $\mathcal{G}$.
(1) For a given $V_{1}$, define $I_{1}=\left\{1_{1}=1,1_{2}, \ldots, 1_{n(1)}\right\} \subset \mathbb{N}$ to be the finite set of positive integers such that $W_{1_{1}}, W_{1_{2}}, \ldots, W_{1_{n(1)}}$ are the only $W_{i}^{\prime} s$ sets for which $V_{1} \cap W_{i} \neq \emptyset$ (if there were infinitely many such $W_{i}^{\prime} s$, then we would have $V_{1} \cap V_{i}^{\prime} s \neq \emptyset$ for infinitely many $i$ 's, which would contradict the star-finiteness of $\left\{V_{i}\right\}_{i \geq 1}$; obviously, we assume that $V_{i} \neq V_{i^{\prime}}$ for $i \neq i^{\prime}$ ). Since $\left\{W_{i}\right\}_{i \geq 1}$ is a cover we have $\bigcup_{i \in I_{1}} W_{i} \cap V_{1}=V_{1}$. (A priori $V_{1}$ can be covered by a proper subfamily of $\left\{V_{2}, \ldots, V_{n}\right\}$ ). Assuming that $i_{1}$ is the greatest number in $I_{1}$ (in particular $i_{1} \geq 1$ ) we set

$$
\varepsilon_{1}=\frac{1}{2} \cdot \frac{1}{2^{i_{1}+1}}>0
$$

We want to apply Lemma 2.27 for the sets

$$
X_{1_{1}}=X_{1}=X_{1} \cap V_{1}, X_{1_{2}} \cap V_{1}, \ldots, X_{1_{n(1)}} \cap V_{1}
$$

which play the role of $X_{0}, \ldots, X_{n}$ in the statement of the Lemma 2.27, and for the sets

$$
W_{1_{1}}=W_{1}=W_{1} \cap V_{1}, W_{1_{2}} \cap V_{1}, \ldots, W_{1_{n(1)}} \cap V_{1}
$$

which play the role of $V_{0}, \ldots, V_{n}$ respectively, and for the positive number $\varepsilon_{1}>0$. Observe that each $X_{i} \cap V_{1}, i \in I_{1}$, has the strong $C^{p}$ extraction property with respect to $V_{1}$. Hence, applying Lemma 2.27, we find a $C^{p}$ diffeomorphism $g_{1}$ of $V_{1} \backslash X_{1}$ onto $V_{1} \backslash\left(X_{1} \backslash U\right)$ which is the identity outside $\left(W_{1} \cap U\right) \cup X_{1}$, carries $\left(X_{l} \backslash X_{1}\right) \cap V_{1}$ into $W_{l} \cap V_{1}$ for each $l>1$ and moves no point more than $\varepsilon_{1}$.
By Remark 2.29 we may also assume that $g_{1}$ sends each set $V_{1, r} \backslash X_{1}$ into $V_{1, r}$. This means in particular that $g_{1}$ refines $\mathcal{G}$. Also, since $g_{1}$ moves no point more than $\varepsilon_{1}$ we cannot have that $x \in V_{i, r}$ and $g_{1}(x) \in V_{i, r^{\prime}}$ for some $i \in I_{1}$ and different $r, r^{\prime} \in \Omega$ (recall that $\operatorname{dist}\left(V_{i, r}, V_{i, r^{\prime}}\right) \geq$ $\frac{1}{2^{i+1}}>\varepsilon_{1}$ ).
Since $W_{1} \cap U \subseteq \overline{W_{1}} \cap U \subseteq V_{1} \cap U$, by making $g_{1}$ be the identity outside $V_{1} \backslash X_{1}$, there exists a well-defined natural extension of $g_{1}$ from $V_{1} \backslash X_{1}$ to $E \backslash X_{1}$. Now we have a $C^{p}$ diffeomorphism $g_{1}$ such that
(a) $g_{1}$ acts from $E \backslash X_{1}$ onto $E \backslash\left(X_{1} \backslash U\right)$.
(b) $g_{1}$ is the identity on $E \backslash\left[\left(W_{1} \cap U\right) \cup X_{1}\right]$. In particular $g_{1}$ is the identity on $(E \backslash U) \backslash X_{1}$.
(c) $g_{1}$ carries $\left(X_{l} \backslash X_{1}\right) \cap V_{1}$ into $W_{l} \cap V_{1}$ for each $l>1$.
(d) We require that if $X_{1}=\emptyset$ or $X_{1, r}=\emptyset$ then $g_{1}$ is the identity on $V_{1}$ or $V_{1, r}$, respectively.
(e) $g_{1}$ moves no point more than $\varepsilon_{1}=\frac{1}{2} \cdot \frac{1}{2^{i}+1}$.
(f) $g_{1}$ sends each set $V_{1, r} \backslash X_{1}$ into $V_{1, r}$ for every $r \in \Omega$, so $g_{1}$ refines $\mathcal{G}$.
(2) Consider now the set $V_{2}$ and define $I_{2}=\left\{2_{1}=2,2_{2}, \ldots, 2_{n(2)}\right\} \subset \mathbb{N}$ to be the finite set of natural numbers such that $W_{2_{1}}, W_{2_{2}}, \ldots, W_{2_{n(2)}}$ are the only $W_{i}^{\prime} s$ sets for which $V_{2} \cap W_{i} \neq \emptyset$ (if there were infinitely many such $W_{i}^{\prime} s$ then $V_{2} \cap V_{i}^{\prime} s \neq \emptyset$ for infinitely many $i$ 's, which would contradict the star-finiteness of $\left\{V_{i}\right\}_{i \geq 1}$ ). Assume that $i_{2}$ is the greatest number in $I_{2}$ (in particular $i_{2} \geq 2$ ), and set

$$
\varepsilon_{2}=\frac{1}{2^{2}} \cdot \frac{1}{2^{i_{2}+1}}>0
$$

Since $\left\{W_{i}\right\}_{i \geq 1}$ is a cover we have

$$
\bigcup_{i \in I_{2}} g_{1}\left(V_{2} \backslash X_{1}\right) \cap W_{i} \cap V_{2}=g_{1}\left(V_{2} \backslash X_{1}\right) \cap V_{2}
$$

Again, we want to apply Lemma 2.26 for the sets

$$
\left\{g_{1}\left(\left(X_{i} \backslash X_{1}\right) \cap V_{2}\right) \cap V_{2}: i \in I_{2}\right\}
$$

playing the role of $X_{0}, \ldots, X_{n}$ in the statement of the lemma, for

$$
\left\{g_{1}\left(V_{2} \backslash X_{1}\right) \cap W_{i} \cap V_{2}: i \in I_{2}\right\}
$$

playing the role of the sets $V_{0} \ldots, V_{n}$ respectively. Here we should recall that $g_{1}\left(X_{i} \backslash X_{1}\right) \subseteq W_{i}$. Observe that by Lemma $2.22(3)$, each $g_{1}\left(\left(X_{i} \backslash X_{1}\right) \cap V_{2}\right) \cap V_{2}$ has the strong $C^{p}$ extraction property with respect to the open set $g_{1}\left(V_{2} \backslash X_{1}\right) \cap V_{2}$. Applying the weak version of Lemma 2.27 to these sets we get a $C^{p}$ diffeomorphism $g_{2}$ of

$$
\left[g_{1}\left(V_{2} \backslash X_{1}\right) \cap V_{2}\right] \backslash\left[g_{1}\left(X_{2} \backslash X_{1}\right) \cap V_{2}\right]=g_{1}\left(V_{2} \backslash\left(X_{1} \cup X_{2}\right)\right) \cap V_{2}
$$

onto

$$
\left[g_{1}\left(V_{2} \backslash X_{1}\right) \cap V_{2}\right] \backslash\left[\left(g_{1}\left(X_{2} \backslash X_{1}\right) \cap V_{2}\right) \backslash U\right]
$$

which is the identity outside

$$
\left(g_{1}\left(V_{2} \backslash\left(X_{1} \cup X_{2}\right)\right) \cap W_{2} \cap V_{2} \cap U\right) \cup\left(g_{1}\left(X_{2} \backslash X_{1}\right) \cap V_{2}\right)
$$

and carries

$$
g_{1}\left(X_{k} \backslash\left(X_{1} \cup X_{2}\right) \cap V_{2}\right) \cap V_{2}
$$

into

$$
g_{1}\left(V_{2} \backslash X_{1}\right) \cap W_{k} \cap V_{2}
$$

for each $k>2$. Moreover using Remark 2.28 one can also assume that $g_{2} \circ g_{1}$ moves no point more than $\varepsilon_{1}+\varepsilon_{2}$.
Because $\overline{g_{1}\left(V_{2} \backslash\left(X_{1} \cup X_{2}\right)\right) \cap W_{2}} \cap U \subseteq g_{1}\left(V_{2} \backslash\left(X_{1} \cup X_{2}\right)\right) \cap V_{2} \cap U$, by letting $g_{2}$ be the identity outside $g_{1}\left(V_{2} \backslash\left(X_{1} \cup X_{2}\right)\right) \cap V_{2}$ there exists a well-defined natural extension of $g_{2}$ to $g_{1}\left(E \backslash\left(X_{1} \cup X_{2}\right)\right)$. To sum up we have the following properties:
(a) $g_{2}$ acts from

$$
g_{1}\left(E \backslash\left(X_{1} \cup X_{2}\right)\right)
$$

onto

$$
\begin{aligned}
g_{1}\left(E \backslash X_{1}\right) \backslash\left[g_{1}\left(X_{2} \backslash X_{1}\right) \backslash U\right] & =E \backslash\left(X_{1} \backslash U\right) \backslash\left[\left(g_{1}\left(\left(X_{2} \backslash X_{1}\right) \cap U\right) \cup g_{1}\left(\left(X_{2} \backslash X_{1}\right)\right) \backslash U\right) \backslash U\right]= \\
& =E \backslash\left(X_{1} \backslash U\right) \backslash\left[\left(X_{2} \backslash X_{1}\right) \backslash U\right]=E \backslash\left(\left(X_{1} \cup X_{2}\right) \backslash U\right)
\end{aligned}
$$

(Here we are using the fact that $g_{1}\left(\left(X_{2} \backslash X_{1}\right) \cap U\right) \subseteq U$ and that $g_{1}$ is the identity outside $U)$.
(b) $g_{2}$ is the identity on $E \backslash\left[\left(g_{1}\left(W_{2} \backslash\left(X_{1} \cup X_{2}\right)\right) \cap W_{2} \cap U\right) \cup\left(g_{1}\left(X_{2} \backslash X_{1}\right) \cap V_{2}\right)\right]$. Since

$$
\left(g_{1}\left(W_{2} \backslash\left(X_{1} \cup X_{2}\right)\right) \cap W_{2} \cap U\right) \cup\left(g_{1}\left(X_{2} \backslash X_{1}\right) \cap V_{2}\right) \subseteq U \cup X_{1} \cup X_{2}
$$

in particular $g_{2}$ is the identity on $E \backslash\left(U \cup\left(X_{1} \cup X_{2}\right)\right)$. Because $g_{1}$ is the identity outside $U$ then $g_{2}$ is the identity on $g_{1}\left(E \backslash\left(U \cup\left(X_{1} \cup X_{2}\right)\right)=g_{1}\left((E \backslash U) \backslash\left(X_{1} \cup X_{2}\right)\right)\right.$.
(c) $g_{2}$ carries

$$
g_{1}\left(X_{l} \backslash\left(X_{1} \cup X_{2}\right)\right) \cap V_{2}
$$

into

$$
W_{l} \cap V_{2}
$$

for each $l>2$.
(d) If $X_{2, r}=\emptyset$ then we require that $g_{2}$ is the identity on $g_{1}\left(V_{2, r} \backslash X_{1}\right) \cap V_{2}$ and in $g_{1}\left(V_{2} \backslash\right.$ $\left.X_{1}\right) \cap V_{2, r}$.
(e) We have

$$
\begin{cases}\left\|g_{2}\left(g_{1}(x)\right)-x\right\| \leq \varepsilon_{1}+\varepsilon_{2}=\frac{1}{2} \cdot \frac{1}{2^{i}+1}+\frac{1}{4} \cdot \frac{1}{2^{i_{2}+1}} & \text { for every } x \in E \backslash\left(X_{1} \cup X_{2}\right)  \tag{2.10}\\ \left\|g_{2}\left(g_{1}(x)\right)-x\right\|<\frac{1}{2^{2+1}} & \text { for every } x \in V_{2} \backslash\left(X_{1} \cup X_{2}\right) .\end{cases}
$$

The first inequality is clear. For the second one, when $x \notin V_{2} \backslash V_{1}$ we have $g_{1}(x)=x$, hence

$$
\left\|g_{2}\left(g_{1}(x)\right)-x\right\|=\left\|g_{2}(x)-x\right\| \leq \varepsilon_{2}=\frac{1}{4} \cdot \frac{1}{2^{i_{2}+1}}<\frac{1}{2^{2+1}}
$$

and if $x \in V_{1} \cap V_{2}$ then $V_{1} \cap V_{2} \neq \emptyset$, so $i_{1} \geq 2$ and

$$
\left\|g_{2}\left(g_{1}(x)\right)-x\right\| \leq \frac{1}{2} \cdot \frac{1}{2^{i_{1}+1}}+\frac{1}{4} \cdot \frac{1}{2^{i_{2}+1}}<\left(\frac{1}{2}+\frac{1}{4}\right) \cdot \frac{1}{2^{2+1}} .
$$

(f) The composition $g_{2} \circ g_{1}$ is a mapping from $E \backslash\left(X_{1} \cup X_{2}\right)$ onto $E \backslash\left(\left(X_{1} \cup X_{2}\right) \backslash U\right)$ and it is the identity outside $(E \backslash U) \backslash\left(X_{1} \cup X_{2}\right)$. Let us check that it refines $\mathcal{G}$. Take $x \in E \backslash\left(X_{1} \cup X_{2}\right)$. If $g_{2}\left(g_{1}(x)\right)=g_{1}(x)$, since $g_{1}$ refines $\mathcal{G}$ we are done. Otherwise we have $g_{2}\left(g_{1}(x)\right) \neq g_{1}(x)$, and then $g_{1}(x), g_{2}\left(g_{1}(x)\right) \in g_{1}\left(V_{2} \backslash\left(X_{1} \cup X_{2}\right)\right) \cap V_{2}$. We have that $x, g_{2}\left(g_{1}(x)\right) \in V_{2}$. We must have $x, g_{2}\left(g_{1}(x)\right) \in V_{2, r}$ for some $r \in \Omega$ (as otherwise $x \in V_{2, r}$ and $g_{2}\left(g_{1}(x)\right) \in V_{2, r^{\prime}}$ a contradiction with (2.10) since $\left.\operatorname{dist}\left(V_{2, r}, V_{2, r^{\prime}}\right) \geq \frac{1}{2^{2+1}}\right)$.
(3) We go on doing this process by successive applications of Lemma [2.27, We want to apply induction to prove that for each $j \geq 3$ we can find a $C^{p}$ diffeomorphism $g_{j}$ such that
(a) $g_{j}$ acts from

$$
g_{j-1} \circ \cdots \circ g_{1}\left(E \backslash \bigcup_{k \leq j} X_{k}\right)
$$

onto

$$
E \backslash\left(\bigcup_{k \leq j} X_{k} \backslash U\right)
$$

(b) $g_{j}$ is the identity on $g_{j-1} \circ \cdots \circ g_{1}\left((E \backslash U) \backslash\left(\bigcup_{k \leq j} X_{k}\right)\right)$.
(c) $g_{j}$ carries

$$
g_{j-1} \circ \cdots \circ g_{1}\left(X_{l} \backslash \bigcup_{k \leq j} X_{k}\right) \cap V_{j}
$$

into

$$
W_{l} \cap V_{j}
$$

for eack $l>j$.
(d) If $X_{j, r}=\emptyset$ then $g_{j}$ is the identity on $g_{j-1} \circ \cdots \circ g_{1}\left(V_{j, r} \backslash \bigcup_{k<j} X_{k}\right) \cap V_{j}$ and on $g_{j-1} \circ \cdots \circ$ $g_{1}\left(V_{j} \backslash \bigcup_{k<j} X_{k}\right) \cap V_{j, r}$.
(e) We have

$$
\begin{cases}\left\|g_{j} \circ \cdots \circ g_{1}(x)-x\right\| \leq \sum_{k=1}^{j} \frac{1}{2^{k}} \cdot \frac{1}{2^{\frac{1}{k}+1}} & \text { for every } x \in E \backslash\left(\bigcup_{k \leq j} X_{k}\right) \\ \left\|g_{j} \circ \cdots \circ g_{1}(x)-x\right\|<\frac{1}{2^{j+1}} & \text { for every } x \in V_{j} \backslash\left(\bigcup_{k \leq j} X_{k}\right),\end{cases}
$$

where $i_{k}$ is the greatest number such that $V_{i_{k}} \cap V_{k} \neq \emptyset$. Let us call for every $k=1, \ldots, j$,

$$
\varepsilon_{k}=\frac{1}{2^{k}} \cdot \frac{1}{2^{i_{k}+1}}>0 .
$$

(f) $g_{j} \circ \cdots \circ g_{1}$ refines $\mathcal{G}$.

Suppose this is true for $j-1$ and let us check this is so for $j$.
The idea is the same as in steps (1) and (2). We first find the set $I_{j}=\left\{j_{1}=j, j_{2}, \ldots, j_{n(j)}\right\} \subset$ $\mathbb{N}$ such that $W_{j_{1}}, W_{j_{2}}, \ldots, W_{j_{n(j)}}$ are the only $W_{i}^{\prime} s$ sets for which $V_{j} \cap W_{i} \neq \emptyset$. Assume $i_{j}$ is the greatest number in $I_{j}$ (in particular $i_{j} \geq j$ ) and set

$$
\varepsilon_{j}=\frac{1}{2^{j}} \cdot \frac{1}{2^{i_{j}+1}}>0
$$

We have that

$$
\bigcup_{i \in I_{j}} g_{j-1} \circ \cdots \circ g_{1}\left(V_{j} \backslash \bigcup_{k<j} X_{k}\right) \cap W_{i} \cap V_{j}=g_{j-1} \circ \cdots \circ g_{1}\left(V_{j} \backslash \bigcup_{k<j} X_{k}\right) \cap V_{j} .
$$

We want to apply Lemma 2.26 for the sets

$$
\left\{g_{j-1} \circ \cdots \circ g_{1}\left(\left(X_{i} \backslash \bigcup_{k<j} X_{k}\right) \cap V_{j}\right) \cap V_{j}: i \in I_{j}\right\}
$$

playing the role of $X_{0}, \ldots, X_{n}$ in the statement of the lemma, for

$$
\left\{g_{j-1} \circ \cdots \circ g_{1}\left(V_{j} \backslash \bigcup_{k<j} X_{k}\right) \cap W_{i} \cap V_{j}: i \in I_{j}\right\}
$$

playing the role of the sets $V_{0} \ldots, V_{n}$ respectively.
Observe that by Lemma 2.22 (3), each $g_{j-1} \circ \cdots \circ g_{1}\left(\left(X_{i} \backslash \bigcup_{k<j} X_{k}\right) \cap V_{j}\right) \cap V_{j}$ has the strong $C^{p}$ extraction property with respect to the open set $g_{j-1} \circ \cdots \circ g_{1}\left(V_{j} \backslash \bigcup_{k<j} X_{k}\right) \cap V_{j}$. Applying the weak version of Lemma 2.27) to these sets we get a $C^{p}$ diffeomorphism $g_{j}$ from

$$
\left[g_{i-1} \circ \cdots \circ g_{1}\left(V_{i} \backslash \bigcup_{j<i} X_{j}\right) \cap V_{i}\right] \backslash\left[g_{j-1} \circ \cdots \circ g_{1}\left(X_{j} \backslash \bigcup_{k<j} X_{k}\right) \cap V_{j}\right]=g_{j-1} \circ \cdots \circ g_{1}\left(V_{j} \backslash \bigcup_{k \leq j} X_{k}\right) \cap V_{j}
$$

onto

$$
\left[g_{j-1} \circ \cdots \circ g_{1}\left(V_{j} \backslash \bigcup_{k<j} X_{k}\right) \cap V_{j}\right] \backslash\left[g_{j-1} \circ \cdots \circ g_{1}\left(X_{j} \backslash \bigcup_{k<j} X_{k}\right) \cap V_{j} \backslash U\right]
$$

which is the identity outside

$$
\left(\left(g_{j-1} \circ \cdots \circ g_{1}\left(V_{j} \backslash \bigcup_{k \leq j} X_{k}\right) \cap W_{j} \cap U\right) \cup\left(g_{j-1} \circ \cdots \circ g_{1}\left(X_{j} \backslash \bigcup_{k<j} X_{k}\right) \cap V_{j}\right)\right.
$$

and carries

$$
g_{j-1} \circ \cdots \circ g_{1}\left(X_{l} \backslash \bigcup_{k \leq j} X_{k}\right) \cap V_{j}
$$

into

$$
W_{l} \cap V_{j}
$$

for each $l>j$. This last property establishes (c).
Define $g_{j}$ to be the natural extension of $g_{j-1} \circ \cdots \circ g_{1}\left(V_{j} \backslash \bigcup_{k \leq j} X_{k}\right) \cap V_{j}$ to $g_{j-1} \circ \cdots \circ g_{1}(E \backslash$
$\bigcup_{k \leq j} X_{k}$ ), so now $g_{j}$ is defined from

$$
\left.g_{j-1} \circ \cdots \circ g_{1}\left(E \backslash \bigcup_{k \leq j} X_{k}\right)\right)=
$$

onto

$$
\begin{aligned}
& {\left[g_{j-1} \circ \cdots \circ g_{1}\left(E \backslash \bigcup_{k<j} X_{k}\right)\right] \backslash\left[g_{j-1} \circ \cdots \circ g_{1}\left(X_{j} \backslash \bigcup_{k<j} X_{k}\right) \backslash U\right]=} \\
& E \backslash\left(\bigcup_{k<j} X_{j} \backslash U\right) \backslash\left[\left(g_{j-1} \circ \cdots \circ g_{1}\left(\left(X_{j} \backslash \bigcup_{k<j} X_{k}\right) \cap U\right) \cup g_{j-1} \circ \cdots \circ g_{1}\left(\left(X_{j} \backslash \bigcup_{k<j} X_{k}\right) \backslash U\right)\right) \backslash U\right]= \\
& =E \backslash\left(\bigcup_{k<j} X_{k} \backslash U\right) \backslash\left[\left(X_{j} \backslash \bigcup_{k<j} X_{k}\right) \backslash U\right]=E \backslash\left(\bigcup_{k \leq j} X_{k} \backslash U\right) .
\end{aligned}
$$

Here we are using that $g_{j-1} \circ \cdots \circ g_{1}\left(\left(X_{j} \backslash \bigcup_{k<j} X_{k}\right) \cap U \subseteq U\right.$ and also (a) from the induction hypothesis. This establishes (a).
We also have that $g_{j}$ is the identity on $(E \backslash U) \backslash\left(\bigcup_{k \leq j} X_{k}\right)=g_{j-1} \circ \cdots \circ g_{1}\left((E \backslash U) \backslash\left(\bigcup_{k \leq j} X_{k}\right)\right)$, which establishes (b).
If $X_{j, r}=\emptyset$ then we let $g_{j}$ be the identity on $g_{j-1} \circ \cdots \circ g_{1}\left(V_{j, r} \backslash \bigcup_{k<j} X_{k}\right) \cap V_{j}$ and on $g_{j-1} \circ \cdots \circ g_{1}\left(V_{j} \backslash \bigcup_{k<j} X_{k}\right) \cap V_{j, r}$. This last property implies (d).
For the property (e), using Remark [2.28 one can assume that $g_{j} \circ \cdots \circ g_{1}$ moves no point more than $\sum_{k=1}^{j} \varepsilon_{k}$, because by the induction hypothesis $g_{j-1} \circ \cdots g_{1}$ moves no point more than $\sum_{k=1}^{j-1} \varepsilon_{k}$. Let us check then that the second property in (2.11) is satisfied. Take $x \in$ $V_{j} \backslash\left(\bigcup_{k \leq j} X_{j}\right)$ and observe that if $g_{k} \circ \cdots \circ g_{1}(x) \neq g_{k-1} \circ \cdots \circ g_{1}(x)$ for some $k=1, \cdots, j$ means that $x \in V_{k}$, hence $k \in I_{j}$ and $i_{k} \geq j$. We can write

$$
\left\|g_{j} \circ \cdots \circ g_{1}(x)-x\right\| \leq \sum_{k \in I_{j}, k \leq j}\left(\frac{1}{2^{k}} \cdot \frac{1}{2^{i_{k}+1}}\right) \leq \sum_{k=1}^{j} \frac{1}{2^{k}} \cdot \frac{1}{2^{j+1}}<\frac{1}{2^{j+1}}
$$

Now, using our induction hypothesis (f) that $g_{j-1} \circ \cdots \circ g_{1}$ refines $\mathcal{G}$, we can prove that $g_{j} \circ \cdots \circ g_{1}$ still refines $\mathcal{G}$. If we take an $x \in E \backslash\left(\bigcup_{k \leq j} X_{k}\right)$ and $g_{j} \circ \cdots \circ g_{1}(x)=g_{j-1} \circ \cdots \circ g_{1}(x)$, since $g_{j-1} \circ \cdots g_{1}$ refines $\mathcal{G}$ we are done. Otherwise $g_{j-1} \circ \cdots g_{1}(x), g_{j} \circ \cdots g_{1}(x) \in g_{j-1} \circ \cdots \circ g_{1}\left(V_{j} \backslash\right.$ $\left.\bigcup_{k<j} X_{k}\right) \cap V_{j}$, so $x, g_{j} \circ \cdots \circ g_{1}(x) \in V_{j}$. We must have $x, g_{j} \circ \cdots \circ g_{1}(x) \in V_{j, r}$ for some $r \in \Omega$. Indeed, otherwise $x \in V_{j, r}$ and $g_{j} \circ \cdots \circ g_{1}(x) \in V_{j, r^{\prime}}$ for different $r, r^{\prime} \in \Omega$, a contradiction with (2.11) $\operatorname{since} \operatorname{dist}\left(V_{j, r}, V_{j, r^{\prime}}\right) \geq \frac{1}{2^{j+1}}$.

Hence we have finished our induction process.
To conclude, note that $g_{j} \circ \cdots \circ g_{1}(x) \neq g_{j-1} \circ \cdots \circ g_{1}(x)$ implies $x \in V_{j}$. This fact ensures the existence of a well-defined $C^{p}$ diffeomorphism $g(x)=\lim _{j \rightarrow \infty} g_{j} \circ \cdots \circ g_{1}(x)$ from $E \backslash X$ onto $E \backslash(X \backslash U)$. The mapping $g$ is the identity on $(E \backslash U) \backslash X$ and because $\left\{V_{j, r}\right\}_{j \geq 1, r \in \Omega}$ is a refinement of $\mathcal{G}, g$ is limited by $\mathcal{G}$.

## 3. Proof of Theorem 1.6

First of all notice that since the result we want to establish is invariant by diffeomorphisms, it is enough to prove it for a $C^{1}$ manifold $M$ diffeomorphic to $E$ in place of $E$. It will be very convenient for us to do so with $M=S^{+}$, the upper sphere of $E \times \mathbb{R}$.
Let $\|\cdot\|$ denote an equivalent norm in $E$ which is LUR and $C^{1}$ (we will also denote by $\|\cdot\|$ the norm of $F$; this will do not do any harm because there will be no risk of confusion). Since $E^{*}$ is separable there always exists such a norm; see [26, Corollary II.4.3] for instance.
Let us define $Y=E \times \mathbb{R}$, with norm

$$
|(u, t)|=\left(\|u\|^{2}+t^{2}\right)^{1 / 2}
$$

and let us denote the upper sphere of $Y$ by

$$
S^{+}:=\left\{(u, t): u \in E, t>0,\|u\|^{2}+t^{2}=1\right\} .
$$

Observe that $S^{+}$is the graph of

$$
s(u)=\sqrt{1-\|u\|^{2}},
$$

which is a $C^{1}$ function ${ }^{2}$ defined on the open unit ball $B_{E}$ of $E$; hence,

$$
d(u)=(u, s(u))=\left(u, \sqrt{1-\|u\|^{2}}\right),
$$

$u \in B_{E}$, defines a $C^{1}$ diffeomorphism of $B_{E}$ onto $S^{+}$. Since $B_{E}$ is obviously $C^{1}$ diffeomorphic with $E$, the upper sphere $S^{+}$, which is a $C^{1}$ submanifold of codimension 1 of $Y$, is also diffeomorphic to $E$. Therefore it will be enough to prove that every continuous mapping $f: S^{+} \rightarrow F$ can be $\varepsilon$-approximated by a mapping $\varphi: S^{+} \rightarrow F$ which is of class $C^{1}$ and has no critical points. As explained in the introduction, this will be done in three steps, the first of which consists in finding a smooth approximation of $f$ whose critical set is a set that we can extract with the help of Theorem 1.4.
In order to find such a smooth approximation, as in [7] we will have to use a partition of unity $\left\{\psi_{n}: n \in \mathbb{N}\right\}$ in $S^{+}$made out of slices of the unit ball of $Y$ by linear functionals $g_{k} \in Y^{*}$, so that the derivative at $y \in S^{+}$of a local sum $\sum \psi_{k}$ will belong to the span of the restrictions to $T_{y} S^{+}$(the tangent space to $S^{+}$at $y$ ) of a finite collection of $g_{k}$. However, the construction of the partition of unity, the technical properties that we will require, and the use that we will make of it, will be much simpler than in that paper.
To construct our partition of unity $\left\{\psi_{n}\right\}_{n \in \mathbb{N}}$ we next translate an old standard argument (going back to Eells in the Hilbert space case, and probably first appearing in [48, p. 28-30], later generalized by Bonic and Frampton [15] for separable Banach spaces with smooth bump functions; we follow [33, Theorem 8.25]) to the upper sphere, as in [65, 36, 7].
Let us denote $S:=S_{|\cdot|}$, the unit sphere of $(Y,|\cdot|)$, and $S^{*}:=S_{|\cdot|^{*}}$, the unit sphere of $\left(Y^{*},|\cdot|^{*}\right)$. The duality mapping of the norm $|\cdot|$, defined as

$$
\begin{aligned}
& D: S \longrightarrow S^{*} \\
& D(x)=|\cdot|^{\prime}(x),
\end{aligned}
$$

is $|\cdot|-|\cdot|^{*}$ continuous since the norm $|\cdot|$ is of class $C^{1}$.
Since the norm $|\cdot|$ is locally uniformly convex we can find, for every $x \in S^{+}$, open slices $R_{x}=\{y \in$ $\left.S: g_{x}(y)>\delta_{x}\right\} \subset S^{+}$and $P_{x}=\left\{y \in S: g_{x}(y)>\delta_{x}^{2}\right\} \subset S^{+}$, where $g_{x}=D(x) \in Y^{*}, 0<\delta_{x}<1$, and $\left|g_{x}\right|^{*}=1=g_{x}(x)$, so that the oscillation of the functions $f$ and $\varepsilon$ on every $P_{x}$ is less than $\varepsilon(x) / 16$. We also assume, with no loss of generality, that $\operatorname{dist}\left(P_{x}, E \times\{0\}\right)>0$.

[^2]Since $Y$ is separable we can select a countable subfamily of $\left\{R_{x}\right\}_{x \in S^{+}}$which covers $S^{+}$. Let us denote this countable subfamily by $\left\{R_{n}\right\}_{n}$, where $R_{n}=R_{x_{n}}=\left\{y \in S: g_{n}(y)>\delta_{n}\right\}$ and $g_{n}\left(x_{n}\right)=1$. Recall that the oscillation of the functions $f$ and $\varepsilon$ on every $P_{n}=P_{x_{n}}=\left\{y \in S: g_{n}(y)>\delta_{n}^{2}\right\}$ is less than $\varepsilon\left(x_{n}\right) / 16$; this implies that

$$
\frac{15}{16} \varepsilon\left(x_{n}\right) \leq \varepsilon(x) \leq \frac{17}{16} \varepsilon\left(x_{n}\right) \text { and }\|f(x)-f(y)\| \leq \frac{\varepsilon\left(x_{n}\right)}{16}
$$

for every $x, y \in P_{n}$. Note also that $\left\{P_{n}\right\}_{n \in \mathbb{N}}$ is an open cover of $S^{+}$.
For each $k \in \mathbb{N}$, let $\theta_{k}: \mathbb{R} \rightarrow[0,1]$ be a $C^{\infty}$ function such that $\theta_{k}(t)=1$ if and only if $t \geq \delta_{k}$, and $\theta(t)=0$ if and only if $t \leq \delta_{k}^{2}$. Next, for each $k \in \mathbb{N}$ we define $\varphi_{k}: S^{+} \rightarrow[0,1]$ by

$$
\varphi_{k}(x)=\theta_{k}\left(g_{k}(x)\right),
$$

and note that the interior of the support of $\varphi_{k}$, which coincides with $\varphi_{k}^{-1}((0,1])$, is the open slice $P_{k}=\left\{y \in S: g_{k}(x)>\delta_{k}^{2}\right\}$.
Now, for $k=1$, define $h_{1}: S^{+} \rightarrow \mathbb{R}$ by

$$
h_{1}(x)=\varphi_{1}(x) .
$$

Notice that the interior of the support of $h_{1}$ is the open set $U_{1}:=P_{1}$.
For $k \geq 2$ let us define $h_{k}: S^{+} \rightarrow \mathbb{R}$ by

$$
h_{k}(x)=\varphi_{k}(x) \prod_{j<k}\left(1-\varphi_{j}(x)\right),
$$

and notice that the interior of the support of $h_{k}$ is the set

$$
U_{k}:=\left\{y \in S^{+}: g_{k}(y)>\delta_{k}^{2} \text { and } g_{j}(y)<\delta_{j} \text { for all } j<k\right\} .
$$

Claim 3.1. The family $\left\{U_{k}\right\}_{k \in \mathbb{N}}$ is a locally finite open covering of $S^{+}$that refines $\left\{P_{k}\right\}_{k \in \mathbb{N}}$. Therefore the functions

$$
\psi_{n}:=\frac{h_{n}}{\sum_{k=1}^{\infty} h_{k}}, n \in \mathbb{N}
$$

define a $C^{1}$ partition of unity in $S^{+}$subordinate to $\left\{P_{k}\right\}_{k \in \mathbb{N}}$.
Proof. Given $j \in \mathbb{N}$, if $x, y \in\left\{y: g_{j}(y)>\delta_{j}\right\}$ and $k>j$ then $\varphi_{j}(y)=1$, hence $h_{k}(y)=0$. Since $\left\{y: g_{j}(y)>\delta_{j}\right\}$ is a neighborhood of $x$ in $S^{+}$this implies that the family of supports of the $h_{k}$ is locally finite. On the other hand, if $h_{k}(x)>0$ for some $x, k$ then $\varphi_{k}(x)>0$, hence $x \in P_{k}$, and this shows that the family of the open supports of the functions $h_{k}$, which coincides with $\left\{U_{k}\right\}_{k \in \mathbb{N}}$, refines $\left\{P_{k}\right\}_{k \in \mathbb{N}}$. It only remains to prove that the family of the open supports of the functions $h_{k}$ is indeed a cover, that is, for every $x \in S^{+}$there exists some $n_{x}$ such that $h_{n_{x}}(x)>0$. We argue by contradiction: assume we had $h_{k}(x)=0$ for all $k \in \mathbb{N}$, then we can show by induction that $\varphi_{n}(x)=0$ for all $n \in \mathbb{N}$, which implies that $x \notin P_{n}$ for all $n$ and contradicts the fact that $\left\{P_{n}\right\}_{n \in \mathbb{N}}$ covers $S^{+}$. Indeed, for $n=1$ we have $0=h_{1}(x)=\varphi_{1}(x)$. Now suppose that we have $\varphi_{1}(x)=\varphi_{2}(x)=\ldots=\varphi_{n}(x)=0$. Then

$$
0=h_{n+1}(x)=\varphi_{n+1}(x) \prod_{j<n+1}\left(1-\varphi_{j}(x)\right)=\varphi_{n+1}(x)=0,
$$

so it is also true that $\varphi_{n+1}(x)=0$.
We will employ the following remarkable fact.
Claim 3.2. For every $k \in \mathbb{N}$ and every $y \in S^{+}$, we have

$$
h_{k}(y)=0 \Longrightarrow D h_{k}(y)=0 .
$$

Proof. Let us assume that $h_{k}(y)=0$. Suppose first that $\varphi_{k}(y)=0$. Computing the derivative of $h_{k}$ at $y$ we get that

$$
D h_{k}(y)=D \varphi_{k}(y) \prod_{j<k}\left(1-\varphi_{j}(y)\right)+\varphi_{k}(y) D\left(\prod_{j<k}\left(1-\varphi_{j}(y)\right)\right)=D \varphi_{k}(y) \prod_{j<k}\left(1-\varphi_{j}(y)\right) .
$$

But $\varphi_{k}(y)=\theta_{k}\left(g_{k}(y)\right)=0$ implies that $D \varphi_{k}(y)=0$, hence $D h_{k}(y)=0$.
If $\varphi_{k}(y) \neq 0$ we must have $\varphi_{j}(y)=1$ for some $j<k$. But again it follows that $D \varphi_{j}(y)=0$, so one can check that also

$$
D h_{k}(y)=\varphi_{k}(y) D\left(\prod_{j<k}\left(1-\varphi_{j}(y)\right)\right)=0 .
$$

Notice that, according to Claim 3.1, for every $x \in S^{+}$there exist a number $n=n_{x}$ and an open neighborhood $V_{x}$ of $x$ in $S^{+}$such that $\sum_{j \leq n} h_{j}(y)>0$ for every $y \in V_{x}$, and $h_{k}(y)=0$ for every $k>n$. This means that $\psi_{k}(y)=0$ for every $k>n$ and every $y \in V_{x}$. More precisely, $n_{x}$ can be chosen as the first such $j$ so that $x \in\left\{y \in S^{+}: g_{j}(y)>\delta_{j}\right\}=R_{j}$ and $V_{x}=R_{j}$.
Let us also call $m=m_{y}$ the largest $j$ such that $h_{j}(y) \neq 0$. Note that $m_{y}$ is also the largest $j$ for which $\psi_{j}(y) \neq 0$. Thus, for every $y \in V_{x}$, we have

$$
m_{y}=\max \left\{j: y \in \psi_{j}^{-1}((0,1])\right\} \leq n_{x} .
$$

In order to calculate the derivatives of this partition of unity in this neighborhood $V_{x}$ of $x$, let us introduce the functions

$$
H_{k}(y)=\frac{h_{k}(y)}{\sum_{j=1}^{n} h_{j}(y)}, \quad y \in V_{x}, \quad k=1, \ldots, n
$$

which are well defined on $V_{x}$ and in fact can be extended as $C^{1}$ smooth functions to an open subset $\mathcal{S}_{n}$ of $Y$ containing $V_{x}$. Specifically, noting that each $h_{k}$ is well-defined and $C^{\infty}$ smooth on the whole $Y$, one can choose $\mathcal{S}_{n}=\bigcup_{i=1}^{n} h_{i}^{-1}((0,1])$. Slightly abusing notation, we will keep denoting these extensions by $H_{k}$, and we will also think of the functions $h_{j}, \psi_{j}, g_{j}, j \leq n$, as being $C^{1}$ smooth functions defined on this open set $\mathcal{S}_{n}$. Therefore, to calculate the derivative of $\psi_{k}$ on $V_{x}$ for $k=1, \ldots, n$, we only have to calculate the derivative of $H_{k}$ at each $y \in V_{x}$ for $k=1, \ldots, n$ and then restrict it to the tangent spaces $T_{y} S^{+}$. The exact expression for the derivative of the functions $H_{1}, \ldots, H_{n}$ on $V_{x}$ will not be particularly interesting or useful to us. The only thing we need to know is that there are $C^{1}$ smooth functions $\sigma_{k, j}, 1 \leq k, j \leq n$, (actually, $\sigma_{k, j}$ will be $C^{\infty}$ smooth) defined on $\mathcal{S}_{n}$ so that

$$
D H_{k}(y)=\sum_{j=1}^{n} \sigma_{k, j}(y) g_{j}
$$

for $k=1, \ldots, n, y \in V_{x}$, and that, in fact, as an immediate consequence of Claim 3.2 and the definition of $m_{y}$ we have

$$
D H_{k}(y)=\sum_{j=1}^{m_{y}} \sigma_{k, j}(y) g_{j} .
$$

for $k=1, \ldots, m_{y}, y \in V_{x}$.
Note that even though each $H_{k}$ is $C^{\infty}$ smooth on an open subset $\mathcal{S}_{n}$ of $Y$, we cannot say that $\psi_{k}$ is $C^{\infty}$ smooth too, because $S^{+}$has not a $C^{\infty}$ smooth submanifold structure. The functions $\psi_{k}$ are just $C^{1}$ because $S^{+}$is just a $C^{1}$ manifold modeled on $E$. Now, if we want to know what the derivative of
the functions $\psi_{k}, k=1, \ldots, n$, on $V_{x}$ looks like, because $H_{k}=\psi_{k}$ on $V_{x}$ and $V_{x}$ is open in $S^{+}$, we only have to restrict $D H_{k}$ to the tangent spaces $T_{y} S^{+}$for each $y \in V_{x}$. Hence we have

$$
D \psi_{k}(y)(v)=\sum_{j=1}^{n} \sigma_{k, j}(y) g_{j}(v)=\sum_{j=1}^{m_{y}} \sigma_{k, j}(y) g_{j}(v)
$$

for each $v \in T_{y} S^{+}$.
This expression can be somewhat misleading at first sight, because one might think that, for $k=$ $1, \ldots, n, D \psi_{k}(y)$ is just a linear combination of the functionals $g_{1}, \ldots, g_{n}$, and this is not exactly so. It is a linear combination of the restrictions of $g_{1}, \ldots, g_{n}$ to $T_{y} S^{+}$, and therefore for every $y \in V_{x}$ it is a different linear combination of different linear functionals $g_{\left.1\right|_{T_{y} S^{+}}}, \ldots, g_{\left.1\right|_{T_{y} S^{+}}}$, each of them defined on a space depending on $y$.
In order to fully clarify this important point, let us calculate the tangent space $T_{y} S^{+}$at $y=\left(u_{y}, t_{y}\right)$. Since $S^{+}$is the graph of the function $s(u)=\sqrt{1-\|u\|^{2}}$, the most natural representation of $T_{y} S^{+}$is given by

$$
\begin{equation*}
T_{y} S^{+}=\left\{(u, t) \in Y=E \times \mathbb{R}: t=L_{y}(u), u \in E\right\}=\left\{\left(u, L_{y}(u)\right): u \in E\right\} \subset Y \tag{3.1}
\end{equation*}
$$

where $L_{y}$ is the derivative $D s\left(u_{y}\right)$ of the function $s$ evaluated at the point $u_{y}=d^{-1}(y) \in B_{E}$ (recall that $d(u)=(u, s(u)))$; in other words, if $y \neq(0,1)$,

$$
L_{y}(w)=-\frac{\left\|u_{y}\right\|}{\sqrt{1-\left\|u_{y}\right\|^{2}}} D\|\cdot\|\left(u_{y}\right)(w)=-\frac{\sqrt{1-t_{y}^{2}}}{t_{y}} D\|\cdot\|\left(u_{y}\right)(w)
$$

for each $w \in E$. Of course we have $L_{(0,1)}(w)=D s(0)(w)=0$ for each $w \in E$, and $T_{(0,1)} S^{+}=E \times\{0\}$. This is the vectorial tangent hyperplane to $S^{+}$at $y$, as opposed to the affine tangent hyperplane to $S^{+}$, which is just $y+T_{y} S^{+}$. Since derivatives of mappings act on vectorial tangent hyperplanes we may forget the affine hyperplanes $y+T_{y} S^{+}$in what follows.
Now, recall that $g_{j} \in Y^{*}$ and therefore these functionals are of the form

$$
g_{j}(u, t)=g_{j}^{1}(u)+g_{j}^{2} t, \quad(u, t) \in E \times \mathbb{R}=Y
$$

where $g_{j}^{1} \in E^{*}$ and $g_{j}^{2} \in \mathbb{R}$.
Therefore the derivative of $g_{\left.j\right|_{S^{+}}}$at $y \in S^{+}$is given by

$$
\begin{equation*}
D g_{j}(y)\left(u, L_{y}(u)\right)=g_{j}\left(u, L_{y}(u)\right)=g_{j}^{1}(u)+g_{j}^{2} L_{y}(u) \tag{3.2}
\end{equation*}
$$

for every $v=\left(u, L_{y}(u)\right) \in T_{y} S^{+}$. Thus, for every $y \in V_{x}$ and every $k=1, \ldots, n$, we have

$$
D \psi_{k}(y)(v)=D \psi_{k}(y)\left(u, L_{y}(u)\right)=\sum_{j=1}^{n} \sigma_{k, j}(y)\left(g_{j}^{1}(u)+g_{j}^{2} L_{y}(u)\right)=\sum_{j=1}^{m_{y}} \sigma_{k, j}(y)\left(g_{j}^{1}(u)+g_{j}^{2} L_{y}(u)\right)
$$

for every $v=\left(u, L_{y}(u)\right) \in T_{y} S^{+}$.
Finally, let us note that the points $x_{n} \in R_{n}$ satisfy that

$$
\frac{15}{16} \varepsilon\left(x_{n}\right) \leq \varepsilon(y) \leq \frac{17}{16} \varepsilon\left(x_{n}\right) \text { for every } y \in P_{n}, \quad \text { and } \quad \sup _{x, y \in P_{n}}\|f(x)-f(y)\| \leq \frac{\varepsilon\left(x_{n}\right)}{16}
$$

The following lemma summarizes the properties of the partition of unity $\left\{\psi_{n}\right\}_{n \in \mathbb{N}}$ which will be most useful to us.
Lemma 3.3. Given two continuous functions $f: S^{+} \rightarrow F$ and $\varepsilon: S^{+} \rightarrow(0, \infty)$, there exists a collection of norm-one linear functionals $\left\{g_{k}\right\}_{k \in \mathbb{N}} \subset Y^{*}$, an open covering $\left\{P_{n}\right\}_{n \in \mathbb{N}}$ of $S^{+}$, and a $C^{1}$ partition of unity $\left\{\psi_{n}\right\}_{n \in \mathbb{N}}$ on $S^{+}$such that:
(1) $\left\{\psi_{n}\right\}_{n \in \mathbb{N}}$ is subordinate to $\left\{P_{n}\right\}_{n \in \mathbb{N}}$.
(2) For every $x \in S^{+}$there exist a neighborhood $V_{x}$ of $x$ in $S^{+}$and a number $n=n_{x} \in \mathbb{N}$ such that $\psi_{m}=0$ on $V_{x}$ for all $m>n$, and the derivatives of the functions $\psi_{1}, \ldots, \psi_{n}$ on $V_{x}$ are of the form

$$
D \psi_{k}(y)(v)=\sum_{j=1}^{n} \sigma_{k, j}(y) g_{j}(v)=\sum_{j=1}^{m_{y}} \sigma_{k, j}(y) g_{j}(v)
$$

for $v \in T_{y} S^{+}$, the tangent hyperplane to $S^{+}$at $y \in S^{+} \cap V_{x}$, and where $m_{y} \leq n$ is the largest number such that $\psi_{m_{y}}(y) \neq 0$. More precisely, if $L_{y}$ denotes the derivative of the function $s(u)=\sqrt{1-\|u\|^{2}}$ evaluated at the point $d^{-1}(y)$, where $d(u)=(u, s(u))$, we have

$$
D \psi_{k}(y)(v)=D \psi_{k}(y)\left(u, L_{y}(u)\right)=\sum_{j=1}^{n} \sigma_{k, j}(y)\left(g_{j}^{1}(u)+g_{j}^{2} L_{y}(u)\right)=\sum_{j=1}^{m_{y}} \sigma_{k, j}(y)\left(g_{j}^{1}(u)+g_{j}^{2} L_{y}(u)\right)
$$

for every $k=1, \ldots, n$, and for every $v=\left(u, L_{y}(u)\right) \in T_{y} S^{+}$, where the functions $\sigma_{k, j}: V_{x} \rightarrow \mathbb{R}$ are of class $C^{1}$, and $g_{j}^{1} \in E^{*}, g_{j}^{2} \in \mathbb{R}, j=1, \ldots, n$.
(3) For every $n \in \mathbb{N}$ there exist a point $y_{n}:=x_{n} \in P_{n}$ such that

$$
\frac{15}{16} \varepsilon\left(y_{n}\right) \leq \varepsilon(y) \leq \frac{17}{16} \varepsilon\left(y_{n}\right) \quad \text { for every } y \in P_{n}, \quad \text { and } \quad \sup _{x, y \in P_{n}}\|f(x)-f(y)\| \leq \frac{\varepsilon\left(y_{n}\right)}{16} .
$$

Now we are ready to start the construction of our approximating function $\varphi: S^{+} \rightarrow F$, which will be of the form

$$
\varphi(x)=\sum_{n=1}^{\infty}\left(f\left(y_{n}\right)+T_{n}(x)\right) \psi_{n}(x),
$$

where the $y_{n}$ are the points given by condition (4) of the preceding lemma, and the operators $T_{n}$ : $Y \rightarrow F$ will be carefully defined below.

Case 1: Assume that $F$ is infinite-dimensional.
We will have to make repeated use of the following fact, whose proof is elementary and can be left to the interested reader.

Lemma 3.4. If $E$ is a Banach space which is isomorphic to $E \oplus E$ then for every finite-codimensional closed subspace $V$ of $E$, there exists a decomposition

$$
E=E_{1} \oplus E_{2} \oplus E_{3},
$$

with factors $E_{1}, E_{2}, E_{3}$ isomorphic to $E$, such that $E_{1} \oplus E_{2} \subset V$.
We start considering a decomposition

$$
E=E_{1,1} \oplus E_{1,2},
$$

with infinite-dimensional factors isomorphic to $E$, and we define a continuous linear surjection $S_{1}: E \rightarrow$ $F$ such that $S_{1}=0$ on $E_{1,2}$. This can be done by taking a continuous linear surjection $R_{1}: E_{1,1} \rightarrow F$ (which exists because by assumption there exists such an operator from $E$ onto $F$ and $E_{1,1}$ is isomorphic to $E$ ), and setting $S_{1}=R_{1} \circ P_{1,1}$, where $P_{1,1}: E_{1,1} \oplus E_{1,2} \rightarrow E_{1,1}$ is the projection onto the first factor associated to this decomposition of $E$. Next, recall that the linear functionals $g_{n} \in Y^{*}=(E \times \mathbb{R})^{*}$ are of the form

$$
g_{j}(u, t)=g_{j}^{1}(u)+g_{j}^{2} t,
$$

where $g_{j}^{1} \in E^{*}$ and $g_{j}^{2} \in \mathbb{R}$. Of course, $\bigcap_{j=1}^{2} \operatorname{Ker} g_{j}^{1}$ is a finite-codimensional subspace of $E$, so by using Lemma 3.4 (2) with $V=E_{1,2} \cap \bigcap_{j=1}^{2} \operatorname{Ker}\left(g_{j}^{1}\right)$, and with $E_{1,2}$ in place of $E$, we may find a decomposition of the second factor $E_{1,2}$,

$$
E_{1,2}=E_{2,1} \oplus E_{2,2} \oplus E_{2,3},
$$

with factors $E_{2,1}, E_{2,2}$ and $E_{2,3}$ isomorphic to $E$, and

$$
E_{2,1} \oplus E_{2,2} \subset \bigcap_{j=1}^{2} \operatorname{Ker} g_{j}^{1},
$$

and we may easily define a bounded linear operator $S_{2}$ from $E$ onto $F$ such that $S_{2}=0$ on $E_{1,1} \oplus E_{2,2} \oplus$ $E_{2,3}$ (this can be done by taking a surjective operator $R_{2}: E_{2,1} \rightarrow F$ and defining $S_{2}=R_{2} \circ P_{2,1} \circ P_{1,2}$, where $P_{1,2}: E_{1,1} \oplus E_{1,2} \rightarrow E_{1,2}$ and $P_{2,1}: E_{2,1} \oplus E_{2,2} \oplus E_{2,3} \rightarrow E_{2,1}$ are the projections associated to the corresponding decompositions).
We continue this process by induction: assuming that we have already defined decompositions $E=$ $E_{1,1} \oplus E_{1,2}, \quad E_{1,2}=E_{2,1} \oplus E_{2,2} \oplus E_{2,3}, \quad E_{2,3}=E_{3,1} \oplus E_{3,2} \oplus E_{3,3}, \ldots, E_{n-1,3}=E_{n, 1} \oplus E_{n, 2} \oplus E_{n, 3}$, and surjective operators

$$
S_{k}: E=E_{1,1} \oplus\left(E_{2,1} \oplus E_{2,2}\right) \oplus \ldots \oplus\left(E_{k-1,1} \oplus E_{k-1,2}\right) \oplus\left(E_{k, 1} \oplus E_{k, 2} \oplus E_{k, 3}\right) \rightarrow F, \quad k=2, \ldots, n
$$

so that $S_{k}$ is zero on all the factors of this decomposition except $E_{k, 1}$, and

$$
E_{k, 1} \oplus E_{k, 2} \subset \bigcap_{j=1}^{k} \operatorname{Ker}\left(g_{j}^{1}\right),
$$

we again apply Lemma 3.4 (2) to write

$$
E_{n, 3}=E_{n+1,1} \oplus E_{n+1,2} \oplus E_{n+1,3},
$$

with factors isomorphic to $E$ and

$$
E_{n+1,1} \oplus E_{n+1,2} \subset \bigcap_{j=1}^{n+1} \operatorname{Ker}\left(g_{j}^{1}\right)
$$

and we define a continuous linear surjection

$$
S_{n+1}: E=E_{1,1} \oplus\left(E_{2,1} \oplus E_{2,2}\right) \oplus \ldots \oplus\left(E_{n, 1} \oplus E_{n, 2}\right) \oplus\left(E_{n+1,1} \oplus E_{n+1,2} \oplus E_{n+1,3}\right) \rightarrow F,
$$

by setting it equal to 0 on all the factors of this decomposition except $E_{n+1,1}$, which is mapped onto $F$.
Having this collection of surjective operators $S_{n}: E \rightarrow F$ at our disposal, we finally define $T_{n}: Y \rightarrow F$ by

$$
T_{n}(u, t)=\frac{\varepsilon\left(y_{n}\right)}{4\left\|S_{n}\right\|} S_{n}(u),
$$

and $\varphi: S^{+} \rightarrow F$ by

$$
\varphi(x)=\sum_{n=1}^{\infty}\left(f\left(y_{n}\right)+T_{n}(x)\right) \psi_{n}(x) .
$$

It is clear that $\varphi$ is well defined and of class $C^{1}$.
In the rest of the proof we will check that this mapping $\varepsilon$-approximates $f$ on $S^{+}$, and that the set of critical points of $\varphi$ is a set which can be diffeomorphically extracted by using Theorem 1.4. Then the proof will be completed by setting $g=\varphi \circ h$, where $h: S^{+} \rightarrow S^{+} \backslash C_{\varphi}$ is a $C^{1}$ diffeomorphism which is close enough to the identity.

Claim 3.5. The mapping $\varphi$ approximates $f$.
Proof. By condition (3) of Lemma 3.3, we know that the oscillation of $f$ in $P_{n}$ is less that $\varepsilon\left(y_{n}\right) / 16$, and by definition of $T_{n}$, we have $\left\|T_{n}\right\| \leq \varepsilon\left(y_{n}\right) / 4$, hence $\left\|T_{n}(x)\right\| \leq \varepsilon\left(y_{n}\right) / 4$ for all $x \in S^{+}$too, because $\|x\|=1$. Now, if $\psi_{n}(x) \neq 0$, then $x \in P_{n}$ and

$$
\begin{align*}
\left\|f\left(y_{n}\right)+T_{n}(x)-f(x)\right\| & \leq\left\|f\left(y_{n}\right)-f(x)\right\|+\left\|T_{n}(x)\right\|  \tag{3.3}\\
& \leq \varepsilon\left(y_{n}\right) / 16+\varepsilon\left(y_{n}\right) / 4=\frac{5}{16} \varepsilon\left(y_{n}\right)<\frac{15}{32} \varepsilon\left(y_{n}\right) \leq \varepsilon(x) / 2 .
\end{align*}
$$

Therefore

$$
\|\varphi(x)-f(x)\|=\left\|\sum_{n=1}^{\infty}\left(f\left(y_{n}\right)+T_{n}(x)-f(x)\right) \psi_{n}(x)\right\| \leq \sum_{n=1}^{\infty}\left\|f\left(y_{n}\right)+T_{n}(x)-f(x)\right\| \psi_{n}(x) \leq \varepsilon(x) / 2 .
$$

Let us now consider the question as to how big the critical set $C_{\varphi}$ can be. We need to calculate the derivative of our function $\varphi$. To this end we first have to examine the expressions for the derivatives of the operators $T_{n}$ restricted to $S^{+}$. These are simpler than those of the $g_{n}$ 's on $S^{+}$, because of the way the operators $T_{n}: Y \rightarrow F$ have been defined. Indeed, for every $v=\left(u, L_{y}(u)\right) \in T_{y} S^{+}$we have that

$$
\begin{equation*}
T_{n}(v)=T_{n}\left(u, L_{y}(u)\right)=\frac{\varepsilon\left(y_{n}\right)}{4\left\|S_{n}\right\|} S_{n}(u), \tag{3.4}
\end{equation*}
$$

and we have that $D\left(T_{\left.n\right|_{S^{+}}}\right)(y)$ is the restriction of $D T_{n}(y)=T_{n}$ to $T_{y} S^{+}$, that is to say, if $v=$ $\left(u, L_{y}(u)\right) \in T_{y} S^{+}$then

$$
\begin{equation*}
D T_{n}(y)\left(u, L_{y}(u)\right)=\frac{\varepsilon\left(y_{n}\right)}{4\left\|S_{n}\right\|} S_{n}(u) . \tag{3.5}
\end{equation*}
$$

Now, recall that, by condition (2) of Lemma 3.3, for every $x \in S^{+}$there is a neighborhood $V_{x}$ of $x$ in $S^{+}$and a number $n=n_{x} \in \mathbb{N}$ such that

$$
\varphi(y)=\sum_{j=1}^{n}\left(f\left(y_{j}\right)+T_{j}(y)\right) \psi_{j}(y)
$$

for every $y \in V_{x}$. Fix $y \in V_{x}$ and recall that for $m=m_{y}$, the largest number $j$ for which $h_{j}(y) \neq 0$ (or $\left.\psi_{j}(y) \neq 0\right)$, we have

$$
\varphi(y)=\sum_{j=1}^{m}\left(f\left(y_{j}\right)+T_{j}(y)\right) \psi_{j}(y)
$$

By using (3.5) and the expression for $D \psi_{j}(y)$ given in Lemma 3.3, we easily see that

$$
\begin{align*}
D \varphi(y)\left(u, L_{y}(u)\right) & =\sum_{j=1}^{n} \psi_{j}(y) \frac{\varepsilon\left(y_{j}\right)}{4\left\|S_{j}\right\|} S_{j}(u)+\sum_{j=1}^{n}\left(g_{j}^{1}(u)+g_{j}^{2} L_{y}(u)\right) \alpha_{n, j}(y)= \\
& =\sum_{j=1}^{m} \psi_{j}(y) \frac{\varepsilon\left(y_{j}\right)}{4\left\|S_{j}\right\|} S_{j}(u)+\sum_{j=1}^{m}\left(g_{j}^{1}(u)+g_{j}^{2} L_{y}(u)\right) \alpha_{m, j}(y) \tag{3.6}
\end{align*}
$$

for every $\left(u, L_{y}(u)\right) \in T_{y} S^{+}, y \in V_{x}$, where the functions $\alpha_{n, j}: V_{x} \subset S^{+} \rightarrow F$ are of class $C^{1}$ because $\alpha_{n, j}(y)=\sum_{i=1}^{n}\left(f\left(y_{i}\right)+T_{i}(y)\right) \sigma_{i, j}(y)$.
Let us now show that the critical set of $\varphi$ is relatively small.

When $n_{x}=1$ we have $\varphi(y)=f\left(y_{1}\right)+T_{1}(y)$ on $V_{x}$, so $D \varphi(y)$ is the restriction of $T_{1}$ to $T_{y} S^{+}$, and equation (3.4) for $n=1$ implies that this restriction is a surjective operator. So it is clear that $\varphi$ has no critical point on $V_{x}$.

Claim 3.6. If $n_{x} \geq 2$ then $C_{\varphi} \cap V_{x}$ is contained in the set

$$
A_{x}:=\left\{y \in S^{+}: E_{n, 2} \subset \operatorname{Ker} L_{y}\right\}
$$

Recall that $L_{y}=D s\left(u_{y}\right)$, where $s(u)=\sqrt{1-\|u\|^{2}}, d(u)=(u, s(u))$, and $u_{y}=d^{-1}(y)$.
Proof. Let us see that, if $y \in V_{x} \backslash A_{x}$ then $D \varphi(y): T_{y} S^{+} \rightarrow F$ is surjective, that is, for every $w \in F$ there exists $v \in T_{y} S^{+}$such that $D \varphi(y)(v)=w$. Let $m=m_{y}$ be the largest number such that $\psi_{m}(y) \neq 0$. Recall that $m \leq n$. Since the operator

$$
S_{m}: E=E_{1,1}\left(\oplus E_{2,1} \oplus E_{2,2}\right) \oplus \ldots \oplus\left(E_{m-1,1} \oplus E_{m-1,2}\right) \oplus\left(E_{m, 1} \oplus E_{m, 2} \oplus E_{m, 3}\right) \rightarrow F
$$

is surjective and equal to zero on all the factors of this decomposition except $E_{m, 1}$ (which is mapped onto $F$ ), we may find $u_{m, 1} \in E_{m, 1}$ so that

$$
S_{m}\left(u_{m, 1}\right)=4 \varepsilon\left(y_{m}\right)^{-1} \psi_{m}(y)^{-1}\left\|S_{m}\right\| w
$$

Now, since $y \notin A_{x}$ there exists some $e_{n, 2} \in E_{n, 2} \backslash$ Ker $L_{y}$, that is to say, $L_{y}\left(e_{n, 2}\right) \neq 0$, and this implies that, if we put

$$
t_{0}:=-\frac{L_{y}\left(u_{m, 1}\right)}{L_{y}\left(e_{n, 2}\right)},
$$

then the vector

$$
u:=u_{m, 1}+t_{0} e_{n, 2},
$$

satisfies that

$$
L_{y}(u)=0 .
$$

But recall that, for every $k \leq n$, we have $E_{k, 1} \oplus E_{k, 2} \subset \bigcap_{j=1}^{k} \operatorname{Ker}\left(g_{j}^{1}\right)$; in particular, $E_{m, 1} \oplus E_{m, 2} \subset$ $\bigcap_{j=1}^{m} \operatorname{Ker}\left(g_{j}^{1}\right)$ because $m \leq n$. Hence, $g_{j}^{1}(u)=0$ for every $1 \leq j \leq m$. It follows that

$$
\sum_{j=1}^{m}\left(g_{j}^{1}(u)+g_{j}^{2} L_{y}(u)\right) \alpha_{m, j}(y)=0
$$

The rest of the operators $S_{1}, \ldots, S_{m-1}$ are zero on $E_{m, 1} \oplus E_{n, 2}$, so we have

$$
S_{j}(u)=0 \text { for every } j=1, \ldots, m-1
$$

and since $S_{m}$ is zero on $E_{n, 2} \subset E_{m, 3}$ we also have $S_{m}\left(t_{0} e_{n, 2}\right)=0$. Therefore, by combining these equalities with equation (3.6, we obtain that

$$
D \varphi(y)\left(u, L_{y}(u)\right)=w,
$$

and the proof of the claim is complete.

Lemma 3.7. For $x \in S^{+}$with $n:=n_{x} \geq 2$, the set $A_{x}$ of Claim 3.6 is of the form

$$
A_{x}=d\left(G\left(f_{x}\right) \cap B_{E}\right),
$$

where $G\left(f_{x}\right)$ is the graph of a continuous mapping $f_{x}: E_{1,1} \oplus\left(E_{2,1} \oplus E_{2,2}\right) \oplus \cdots \oplus\left(E_{n, 1} \oplus E_{n, 3}\right) \rightarrow E_{n, 2}$.

Proof. Note that, by Claim 3.6,

$$
A_{x}=\left\{d(u)=\left(u, \sqrt{1-\|u\|^{2}}\right):\|u\|<1, u \in \mathcal{A}_{x}\right\}
$$

where

$$
\mathcal{A}_{x}:=\bigcap_{e \in E_{n, 2}}\{u \in E \backslash\{0\}:\langle D\|\cdot\|(u), e\rangle=0\} \cup\{0\}
$$

Let us denote $E_{n, 2}^{\prime}=E_{1,1} \oplus\left(E_{2,1} \oplus E_{2,2}\right) \oplus \cdots \oplus\left(E_{n, 1} \oplus E_{n, 3}\right)$, and let us see that there exists a mapping $f_{x}: E_{n, 2}^{\prime} \rightarrow E_{n, 2}$ such that $\mathcal{A}_{x}=G\left(f_{x}\right)=\left\{w+f_{x}(w): w \in E_{n, 2}^{\prime}\right\}$.
Pick a point $w \in E_{n, 2}^{\prime}$. Note that the function $E_{n, 2} \ni v \mapsto \psi_{w}(v):=\|w+v\|^{2}$ is convex and continuous, and satisfies $\lim _{\|v\| \rightarrow \infty} \psi_{w}(v)=\infty$, hence, since $E_{n, 2}$ is reflexive, $\psi_{w}$ attains a minimum at some point $v_{w} \in E_{n, 2}$; in fact this minimum point $v_{w}$ is unique because the norm $\|\cdot\|$ is strictly convex. Let us denote

$$
f_{x}(w):=v_{w} .
$$

Note that the critical points of $\psi_{w}$, with $w \neq 0$, are exactly the points $v \in E_{n, 2}$ such that

$$
\left.\frac{d}{d t}\|w+v+t e\|^{2}\right|_{t=0}=0 \text { for every } e \in E_{n, 2}
$$

or equivalently

$$
\|w+v\|\langle D\|\cdot\|(w+v), e\rangle=0 \text { for every } e \in E_{n, 2}
$$

which in turn is equivalent to saying that $w+v \in \mathcal{A}_{x}$; we let $f_{x}(0)=v_{0}=0$.
Therefore the unique point $v \in E_{n, 2}$ so that $w+v \in \mathcal{A}_{x}$ is the point $v=f_{x}(w)$. This shows that $\mathcal{A}_{x}$ is the graph of the function $f_{x}$.
Now let us see that the function $f_{x}: E_{n, 2}^{\prime} \rightarrow E_{n, 2}$ is continuous. Suppose $f_{x}$ is discontinuous at $w_{0}$ and let $v_{0}:=f_{x}\left(w_{0}\right)$. Then there exist sequences $w_{k} \rightarrow w_{0}$ in $E_{n, 2}^{\prime}$ and $v_{k}:=f_{x}\left(w_{k}\right)$ in $E_{n, 2}$ and a number $\varepsilon_{0}>0$ so that

$$
\begin{equation*}
\left\|v_{k}-v_{0}\right\| \geq \varepsilon_{0} \text { for all } k \in \mathbb{N} . \tag{3.7}
\end{equation*}
$$

From the previous argument we know that the point $v_{k}$ is characterized as being the unique point $v_{k} \in E_{n, 2}$ for which we have

$$
\begin{equation*}
\left\|w_{k}+v_{k}\right\| \leq\left\|w_{k}+v_{k}+e\right\| \text { for all } e \in E_{n, 2} \tag{3.8}
\end{equation*}
$$

and similarly $v_{0}$ is the unique point $v_{0} \in E_{n, 2}$ for which

$$
\begin{equation*}
\left\|w_{0}+v_{0}\right\| \leq\left\|w_{0}+v_{0}+e\right\| \text { for all } e \in E_{n, 2} \tag{3.9}
\end{equation*}
$$

By taking $e=-v_{k}$ in (3.8) we learn that

$$
\left\|v_{k}\right\|-\left\|w_{k}\right\| \leq\left\|w_{k}+v_{k}\right\| \leq\left\|w_{k}\right\|,
$$

hence $\left\|v_{k}\right\| \leq 2\left\|w_{k}\right\|$, and because $\left\|w_{k}\right\|$ converges to $\left\|w_{0}\right\|$ we deduce that $\left(v_{k}\right)$ is bounded. Since $E_{n, 2}$ is reflexive, this implies that $\left(v_{k}\right)$ has a subsequence that weakly converges to a point $\xi_{0} \in E_{n, 2}$. We keep denoting this subsequence by $\left(v_{k}\right)$.
Now, if we take $e=-v_{k}+e^{\prime}$ in (3.8), with $e^{\prime} \in E_{n, 2}$, we obtain

$$
\left\|w_{k}+v_{k}\right\| \leq\left\|w_{k}+e^{\prime}\right\| \text { for all } e^{\prime} \in E_{n, 2}
$$

This implies (using the facts that $v_{k} \rightharpoonup \xi_{0}$ and $w_{k} \rightarrow w_{0}$, and the weak lower semicontinuity of the norm) that

$$
\begin{equation*}
\left\|w_{0}+\xi_{0}\right\| \leq \liminf _{k \rightarrow \infty}\left\|w_{k}+v_{k}\right\| \leq \liminf _{k \rightarrow \infty}\left\|w_{k}+e^{\prime}\right\|=\left\|w_{0}+e^{\prime}\right\| \text { for all } e^{\prime} \in E_{n, 2} \tag{3.10}
\end{equation*}
$$

That is, we have shown that

$$
\begin{equation*}
\left\|w_{0}+\xi_{0}\right\| \leq\left\|w_{0}+e^{\prime}\right\| \text { for all } e^{\prime} \in E_{n, 2} \tag{3.11}
\end{equation*}
$$

By taking $e^{\prime}=\xi_{0}+\xi$ with $\xi \in E_{n, 2}$ we conclude that

$$
\left\|w_{0}+\xi_{0}\right\| \leq\left\|w_{0}+\xi_{0}+\xi\right\| \text { for all } \xi \in E_{n, 2} .
$$

According to (3.9), $v_{0}$ is the only point which can satisfy this inequality. Hence $\xi_{0}=v_{0}$.
But (3.10) tells us even more: by taking $e^{\prime}=\xi_{0}$ we also learn that there exists a subsequence ( $w_{k_{j}}$ ) of $\left(w_{k}\right)$ such that

$$
\left\|w_{k_{j}}+v_{k_{j}}\right\| \rightarrow\left\|w_{0}+\xi_{0}\right\| .
$$

Since we also know that $w_{k_{j}}+v_{k_{j}}$ converges to $w_{0}+\xi_{0}$ weakly and the norm $\|\cdot\|$ is locally uniformly convex (hence $\|\cdot\|$ has the Kadec-Klee property), this implies that $w_{k_{j}}+v_{k_{j}}$ converges to $w_{0}+\xi_{0}=$ $w_{0}+v_{0}$ in the norm topology as well. As we also have $\lim _{j \rightarrow \infty} w_{k_{j}}=w_{0}$ in norm, we deduce that $\lim _{j \rightarrow \infty}\left\|v_{k_{j}}-v_{0}\right\|=0$, which contradicts (3.7).

Now we can easily finish the proof of Theorem 1.6. By Claim 3.6 and Lemma 3.7, we see that $C_{\varphi}$ is a diffeomorphic image in $S^{+}$of a relatively closed set $Z$ of the open unit ball $B_{E}$ of $E$ which has the property of being locally contained in the graph of a continuous function defined on a complemented subspace of infinite codimension in $E$. Indeed, let $Z:=d^{-1}\left(C_{\varphi}\right) \subset B_{E}$. Since $C_{\varphi}$ is closed in $S^{+}, Z$ is relatively closed in $B_{E}$. Also if we take $z \in Z$ then, according to Lemma 3.7 applied to $x=d(z) \in S^{+}$, $n=n_{x}$, and $V_{x}$, for a neighborhood $U_{z}:=d^{-1}\left(V_{x}\right)$ of $z$, we have $Z \cap U_{z} \subseteq G\left(f_{x}\right)$, where

$$
G\left(f_{x}\right)=\left\{u=(w, v) \in E_{n, 2}^{\prime} \oplus E_{n, 2}=E: v=f_{x}(w)\right\} .
$$

Observe that $E$ has $C^{1}$ smooth partitions of unity since $E$ has a separable dual. Therefore we may apply Theorem 1.4 to find a $C^{1}$ diffeomorphism which extracts $C_{\varphi}$ from $S^{+}$; more precisely, there exists a diffeomorphism $h: S^{+} \rightarrow S^{+} \backslash C_{\varphi}$ which, in addition, is limited by the open cover $\mathcal{G}$ that we next define. Recall that we have

$$
\begin{equation*}
\|\varphi(x)-f(x)\| \leq \varepsilon(x) / 2 \tag{3.12}
\end{equation*}
$$

for all $x \in S^{+}$. Since $\varphi$ and $\varepsilon$ are continuous, for every $z \in S^{+}$there exists $\delta_{z}>0$ so that if $x, y \in B\left(z, \delta_{z}\right)$ then $\|\varphi(y)-\varphi(x)\| \leq \varepsilon(z) / 4 \leq \varepsilon(x) / 2$. We set $\mathcal{G}=\left\{B\left(x, \delta_{x}\right): x \in S^{+}\right\}$.
Finally, let us define

$$
g=\varphi \circ h .
$$

Since $h$ is limited by $\mathcal{G}$ we have that, for any given $x \in S^{+}$, there exists $z \in S^{+}$such that $x, h(x) \in$ $B\left(z, \delta_{z}\right)$, and therefore $\|\varphi(h(x))-\varphi(x)\| \leq \varepsilon(z) / 4$, that is, we have that

$$
\|g(x)-\varphi(x)\| \leq \varepsilon(z) / 4 \leq \varepsilon(x) / 2
$$

By combining this inequality with (3.12), we obtain that

$$
\|g(x)-f(x)\| \leq \varepsilon(x)
$$

for all $x \in S^{+}$. Besides, it is clear that $g$ does not have any critical point: since $h(x) \notin C_{\varphi}$, we have that the linear map $D \varphi(h(x)): T_{h(x)} S^{+} \rightarrow F$ is surjective, and $D h(x): T_{x} S^{+} \rightarrow T_{h(x)} S^{+}$is a linear isomorphism, so $D g(x)=D \varphi(h(x)) \circ D h(x)$ is a linear surjection from $T_{x} S^{+}$onto $F$ for every $x \in S^{+}$.
Case 2: Assume that $F=\mathbb{R}^{m}$. The main idea of the proof is very similar to that of Case 1. The fact that $F$ is finite dimensional will allow us dispense with the hypothesis that $E=E \oplus E$. We will use the same partition of unity $\left\{\psi_{n}\right\}_{n \in \mathbb{N}}$ provided by Lemma 3.3. We will decompose $E$ inductively as follows. Since $\operatorname{Kerg} 1_{1}^{1}$ has infinite dimension we can write

$$
E=E_{1} \oplus G_{1},
$$

where $E_{1}=\mathbb{R}^{m}$ and $G_{1} \subseteq \operatorname{Ker} g_{1}^{1}$. Then $G_{1} \cap \bigcap_{j=1}^{2} \operatorname{Ker} g_{j}^{1}$ has codimension 0 or 1 in $G_{1}$, which is infinite-dimensional, and we can write

$$
E=E_{1} \oplus\left(E_{2,1} \oplus E_{2,2} \oplus G_{2}\right),
$$

where $E_{2,1}=\mathbb{R}^{m}, E_{2,2}=\{0\}$ or $E_{2,2}=\mathbb{R}, G_{1}=E_{2,1} \oplus E_{2,2} \oplus G_{2}$ for some $G_{2}$ with $\operatorname{dim} G_{2}=\infty$, and

$$
E_{2,1} \oplus G_{2}=\bigcap_{j=1}^{2} \operatorname{Ker} g_{j}^{1} \cap G_{1} \subseteq \bigcap_{j=1}^{2} \operatorname{Ker} g_{j}^{1} .
$$

Inductively, we can write

$$
\begin{equation*}
E=E_{1} \oplus\left(E_{2,1} \oplus E_{2,2}\right) \oplus \cdots \oplus\left(E_{n-1,1} \oplus E_{n-1,2}\right) \oplus\left(E_{n, 1} \oplus E_{n, 2} \oplus G_{n}\right), \tag{3.13}
\end{equation*}
$$

where $E_{1}, E_{2,1} \ldots, E_{n, 1}=\mathbb{R}^{m}, E_{2,2}, \ldots, E_{n, 2}$ are subspaces of dimension 0 or 1 ,

$$
E_{n, 2} \oplus G_{n} \subseteq \bigcap_{j=1}^{n} \operatorname{Ker} g_{j}^{1},
$$

and

$$
G_{k}=\left(E_{k+1,1} \oplus E_{k+1,2} \oplus G_{k+1}\right)
$$

for every $k=1, \ldots, n$.
Now, for each $n \in \mathbb{N}$, we define a continuous linear surjection $S_{n}: E \rightarrow F$ by setting it to be 0 on all the factors of the decomposition (3.13) except on $E_{n, 1}$, which is mapped onto $F=\mathbb{R}^{m}$, and we construct our approximating function $\varphi$ exactly as in the proof of Theorem 1.6. At this point, we only need to show the following variant of Claim 3.6 (in which $G_{n}$ replaces the subspace $E_{n, 2}$ of the previous proof).

Claim 3.8. If $n_{x} \geq 2$ then $C_{\varphi} \cap V_{x}$ is contained in the set

$$
A_{x}:=\left\{y \in S^{+}: G_{n} \subset \operatorname{Ker} L_{y}\right\} .
$$

Recall that $L_{y}=D s\left(u_{y}\right)$, where $s(u)=\sqrt{1-\|u\|^{2}}, d(u)=(u, s(u))$, and $u_{y}=d^{-1}(y)$.
Proof. Let us see that, if $y \in V_{x} \backslash A_{x}$ then $D \varphi(y): T_{y} S^{+} \rightarrow F$ is surjective, that is, for every $w \in F$ there exists $v \in T_{y} S^{+}$such that $D \varphi(y)(v)=w$. Let $m=m_{y}$ be the largest number such that $\psi_{m}(y) \neq 0$. Recall that $m \leq n$. Since the operator

$$
S_{m}: E=E_{1,1}\left(\oplus E_{2,1} \oplus E_{2,2}\right) \oplus \ldots \oplus\left(E_{m-1,1} \oplus E_{m-1,2}\right) \oplus\left(E_{m, 1} \oplus E_{m, 2} \oplus G_{m}\right) \rightarrow F
$$

is surjective and equal to zero on all the factors of the decomposition (3.13) except on $E_{m, 1}$ (which is mapped onto $F$ ), we may find $u_{m, 1} \in E_{m, 1}$ so that

$$
S_{m}\left(u_{m, 1}\right)=4 \varepsilon\left(y_{m}\right)^{-1} \psi_{m}(y)^{-1}\left\|S_{m}\right\| w .
$$

Now, since $y \notin A_{x}$ there exists $e_{n} \in G_{n} \backslash \operatorname{Ker} L_{y}$. If we set

$$
t_{0}:=-\frac{L_{y}\left(u_{m, 1}\right)}{L_{y}\left(e_{n}\right)},
$$

then the vector

$$
u:=u_{m, 1}+t_{0} e_{n},
$$

satisfies that

$$
L_{y}(u)=0 .
$$

But recall that, for every $k \leq n$, we have $E_{k, 1} \oplus G_{k} \subset \bigcap_{j=1}^{k} \operatorname{Ker}\left(g_{j}^{1}\right)$; in particular, $E_{m, 1} \oplus G_{m} \subset$ $\bigcap_{j=1}^{m} \operatorname{Ker}\left(g_{j}^{1}\right)$ because $m \leq n$. Hence, $g_{j}^{1}(u)=0$ for every $1 \leq j \leq m$. It follows that

$$
\sum_{j=1}^{m}\left(g_{j}^{1}(u)+g_{j}^{2} L_{y}(u)\right) \alpha_{m, j}(y)=0
$$

The rest of the operators $S_{1}, \ldots, S_{m-1}$ are zero on $E_{m, 1} \oplus G_{n}$, so we have

$$
S_{j}(u)=0 \text { for every } j=1, \ldots, m-1,
$$

and since $S_{m}$ is zero on $G_{n} \subset G_{m}$ we also have $S_{m}\left(t_{0} e_{n}\right)=0$. Therefore, by combining these equalities with equation 3.6, we obtain that

$$
D \varphi(y)\left(u, L_{y}(u)\right)=w
$$

and the proof of the claim is complete.
Then we also have the following.
Lemma 3.9. For $x \in S^{+}$with $n:=n_{x} \geq 2$, the set $A_{x}$ of Claim 3.8 is of the form

$$
A_{x}=d\left(G\left(f_{x}\right) \cap B_{E}\right),
$$

where $G\left(f_{x}\right)$ is the graph of a continuous mapping $f_{x}: E_{1,1} \oplus\left(E_{2,1} \oplus E_{2,2}\right) \oplus \cdots \oplus\left(E_{n, 1} \oplus E_{n, 2}\right) \rightarrow G_{n}$. Proof. Repeat the proof Lemma 3.7, just replacing $E_{n, 2}$ with $G_{n}$.
Let $Z:=d^{-1}\left(C_{\varphi}\right) \subset B_{E}$. According to Lemma 3.9 applied to $x=d(z) \in S^{+}, n=n_{x}$, and $V_{x}$, for a neighborhood $U_{z}:=d^{-1}\left(V_{x}\right)$ of $z$, we have $Z \cap U_{z} \subseteq G\left(f_{x}\right)$, where

$$
G\left(f_{x}\right)=\left\{u=(w, v) \in G_{n}^{\prime} \oplus G_{n}=E: v=f_{x}(w)\right\},
$$

with $G_{n}^{\prime}$ denoting $E_{1} \oplus\left(E_{2,1} \oplus E_{2,2}\right) \oplus \cdots \oplus\left(E_{n-1,1} \oplus E_{n-1,2}\right) \oplus\left(E_{n, 1} \oplus E_{n, 2}\right)$. Since $G_{n}$ is infinitedimensional, we may use Theorem [1.4, and the rest of the proof goes exactly as in Case 1.

Remark 3.10. Observe that in the infinite-dimensional case we could have asked $A_{x}$ to be

$$
A_{x}:=\left\{y \in S^{+}: E_{n, 3} \subset \operatorname{Ker} L_{y}\right\},
$$

using $E_{n, 3}$ instead of $E_{n, 2}$ and requiring that $E_{n, 1} \oplus E_{n, 3} \subseteq \bigcap_{j=1}^{n} \operatorname{Kerg}_{j}^{1}$.

## 4. Proof of Theorem 1.7

First of all let us note that since $E$ admits a $C^{1}$ equivalent norm, $E$ cannot contain a closed subspace isomorphic to $\ell_{1}$. Furthermore, as noted following the statement of Theorem 1.7, condition (2) implies that the basis $\left\{e_{n}\right\}_{n \in \mathbb{N}}$ is unconditional, and therefore by [50, Theorem 1.c.9], is shrinking, that is, we have that $E^{*}=\overline{\operatorname{span}}\left\{e_{n}^{*}: n \in \mathbb{N}\right\}$, where $\left\{e_{n}^{*}\right\}_{n \in \mathbb{N}}$ are the biorthogonal functionals associated to $\left\{e_{n}\right\}_{n \in \mathbb{N}}$ (in particular, $E^{*}$ is separable).
We keep using the notations $Y=E \times \mathbb{R}$ and $S^{+}$from the proof of Theorem 1.6. As in the case of Theorem 1.6, it will be enough to prove Theorem 1.7 with $S^{+}$in place of $E$. We define $e_{0}=(0,1) \in Y$, $e_{0}^{*}: Y=E \times \mathbb{R} \rightarrow \mathbb{R}$ by

$$
e_{0}^{*}(u, t)=t,
$$

and by slightly abusing notation we identify $e_{n} \in E$ to $\left(e_{n}, 0\right) \in Y$ and also extend the $e_{n}^{*} \in E^{*}$ to $e_{n}^{*}: Y=E \times \mathbb{R} \rightarrow \mathbb{R}$ by

$$
e_{n}^{*}(u, t)=e_{n}^{*}(u)=u_{n} \quad \text { for all } \quad u=\sum_{j=1}^{\infty} u_{j} e_{j}, \in E, t \in \mathbb{R}
$$

Then we may we consider $\left\{e_{n}\right\}_{n \in \mathbb{N}} \cup\left\{e_{0}\right\}$ as a basis of $E \times \mathbb{R}=Y$ with associated coordinate functionals $\left\{e_{n}^{*}\right\}_{n \in \mathbb{N}} \cup\left\{e_{0}^{*}\right\}$, and we have that this basis is also shrinking.
Next we are going to construct a partition of unity in $S^{+}$, quite similar but not identical to that of the proof of Theorem 1.6
Since the norm $|\cdot|$ is locally uniformly convex we can find, for every $x \in S^{+}$, open slices $R_{x}=\{y \in$ $\left.S: f_{x}(y)>\delta_{x}\right\} \subset S^{+}$and $P_{x}=\left\{y \in S: f_{x}(y)>\delta_{x}^{4}\right\} \subset S^{+}$, where $f_{x} \in Y^{*}, 0<\delta_{x}<1$, and $\left|f_{x}\right|^{*}=1=f_{x}(x)$, so that the oscillation of the functions $f$ and $\varepsilon$ on every $P_{x}$ is less than $\varepsilon(x) / 16$. We also assume, with no loss of generality, that $\operatorname{dist}\left(P_{x}, E \times\{0\}\right)>0$.
Since $Y$ is separable we can select a countable subfamily of $\left\{R_{x}\right\}_{x \in S^{+}}$, which covers $S^{+}$. Let us denote this countable subfamily by $\left\{R_{n}\right\}_{n}$, where $R_{n}=R_{x_{n}}=\left\{y \in S: f_{n}(y)>\delta_{n}\right\}$ and $f_{n}\left(x_{n}\right)=1$. Recall that the oscillation of the functions $f$ and $\varepsilon$ on every $P_{n}=P_{x_{n}}=\left\{y \in S: f_{n}(y)>\delta_{n}^{4}\right\}$ is less than $\varepsilon\left(x_{n}\right) / 16$, and this implies that

$$
\frac{15}{16} \varepsilon\left(x_{n}\right) \leq \varepsilon(x) \leq \frac{17}{16} \varepsilon\left(x_{n}\right) \text { and }\|f(x)-f(y)\| \leq \frac{\varepsilon\left(x_{n}\right)}{16}
$$

for every $x, y \in P_{n}$. Note that $\left\{P_{n}\right\}_{n \in \mathbb{N}}$ is an open cover of $S^{+}$.

- For $\mathbf{k}=1$, since $\operatorname{span}\left\{e_{n}^{*}: n \in \mathbb{N}\right\}$ is dense in $E^{*}$, we may find numbers $N_{1} \in \mathbb{N}, \epsilon_{1}, \gamma_{1} \in(0,1)$ with $\epsilon_{1}>\gamma_{1}$, and $\beta_{1,0}, \ldots, \beta_{1, N_{1}} \in \mathbb{R}$ with $\beta_{1,0}>0$ so that the functional $g_{1}$ defined by

$$
g_{1}:=\sum_{j=0}^{N_{1}} \beta_{1, j} e_{j}^{*}
$$

has norm 1 and satisfies

$$
\left\{x \in S: f_{1}(x)>\delta_{1}^{2}\right\} \subset\left\{x \in S: g_{1}(x)>\epsilon_{1}\right\} \subset\left\{x \in S: g_{1}(x)>\gamma_{1}\right\} \subset\left\{x \in S: f_{1}(x)>\delta_{1}^{3}\right\}
$$

Let us define

$$
\begin{aligned}
& h_{1}: S^{+} \longrightarrow \mathbb{R} \\
& h_{1}=\theta_{1}\left(g_{1}\right),
\end{aligned}
$$

where $\theta_{1}: \mathbb{R} \rightarrow[0,1]$ is a $C^{\infty}$ function satisfying

$$
\begin{array}{r}
\theta_{1}(t)=0 \text { if and only if } t \leq \gamma_{1} \\
\theta_{1}(t)=1 \text { if and only if } t \geq \epsilon_{1} .
\end{array}
$$

Note that the interior of the support of $h_{1}$ is the open set $U_{1}:=\left\{x \in S^{+}: g_{1}(x)>\gamma_{1}\right\}$.

- For $\mathbf{k}=\mathbf{2}$. We may again use the density of $\operatorname{span}\left\{e_{n}^{*}: n \in \mathbb{N}\right\}$ in $E^{*}$, in order to find numbers $N_{2} \in \mathbb{N}, \gamma_{2}, \epsilon_{2} \in(0,1)$ with $\gamma_{2}<\epsilon_{2}$, and $\beta_{2,0}, \ldots, \beta_{2, N_{2}} \in \mathbb{R}$ so that the linear functional

$$
g_{2}:=\sum_{j=0}^{N_{2}} \beta_{2, j} e_{j}^{*}
$$

has norm 1 and satisfies

$$
\left\{x \in S: f_{2}(x)>\delta_{2}^{2}\right\} \subset\left\{x \in S: g_{2}(x)>\epsilon_{2}\right\} \subset\left\{x \in S: g_{2}(x)>\gamma_{2}\right\} \subset\left\{x \in S: f_{2}(x)>\delta_{2}^{3}\right\}
$$

We may assume without loss of generality that $N_{1} \leq N_{2}$ (otherwise we may set $\beta_{2, j}=0$ for $N_{2}<j \leq$ $N_{1}$ and take a new $N_{2}$ equal to $N_{1}$ ).
Now we define

$$
\begin{aligned}
& h_{2}: S^{+} \longrightarrow \mathbb{R} \\
& h_{2}=\theta_{2}\left(g_{2}\right)\left(1-\theta_{1}\left(g_{1}\right)\right),
\end{aligned}
$$

where $\theta_{2}: \mathbb{R} \rightarrow[0,1]$ is a $C^{\infty}$ function satisfying:

$$
\begin{aligned}
& \theta_{2}(t)=0 \text { if and only if } t \leq \gamma_{2} \\
& \theta_{2}(t)=1 \text { if and only if } t \geq \epsilon_{2}
\end{aligned}
$$

Notice that the interior of the support of $h_{2}$ is the open set

$$
U_{2}=\left\{x \in S^{+}: g_{1}(x)<\epsilon_{1}, g_{2}(x)>\gamma_{2}\right\}
$$

- For $\mathbf{k}=\mathbf{3}$, By density of $\operatorname{span}\left\{e_{n}^{*}: n \in \mathbb{N}\right\}$ in $E^{*}$ we may pick numbers $N_{3} \in \mathbb{N}, \gamma_{3}, \epsilon_{3} \in(0,1)$ with $\epsilon_{3}>\gamma_{3}$, and $\beta_{3,0}, \ldots, \beta_{3, N_{3}} \in \mathbb{R}$ so that, for

$$
g_{3}:=\sum_{j=0}^{N_{3}} \beta_{3, j} e_{j}^{*}
$$

we have that $g_{3} \in S^{*}$ and

$$
\left\{x \in S: f_{3}(x)>\delta_{3}^{2}\right\} \subset\left\{x \in S: g_{3}(x)>\epsilon_{3}\right\} \subset\left\{x \in S: g_{3}(x)>\gamma_{3}\right\} \subset\left\{x \in S: f_{3}(x)>\delta_{3}^{3}\right\}
$$

Again we may assume without loss of generality that $N_{2} \leq N_{3}$.
We define

$$
\begin{aligned}
& h_{3}: S^{+} \longrightarrow \mathbb{R} \\
& h_{3}=\theta_{3}\left(g_{3}\right) \prod_{j=1}^{2}\left(1-\theta_{j}\left(g_{j}\right)\right)
\end{aligned}
$$

where $\theta_{3}: \mathbb{R} \rightarrow[0,1]$ is a $C^{\infty}$ function satisfying

$$
\begin{array}{cl}
\theta_{3}(t)=0 & \text { if and only if } t \leq \gamma_{3} \\
\theta_{3}(1)=1 & \text { if and only if } t \geq \epsilon_{3}
\end{array}
$$

Clearly the interior of the support of $h_{3}$ is the set

$$
U_{3}=\left\{x \in S^{+}: g_{1}(x)<\epsilon_{1}, g_{2}(x)<\epsilon_{2} \text { and } g_{3}(x)>\gamma_{3}\right\}
$$

We continue this process by induction.

- Assume that, in the steps $j=2, \ldots, k$, with $k \geq 2$, we have selected points $y_{j} \in S^{+}$, positive integers $N_{1} \leq N_{2} \leq \ldots \leq N_{k}$, and constants $\gamma_{j}, \epsilon_{j} \in(0,1), \beta_{j, i} \in \mathbb{R}$ so that the functionals

$$
g_{j}:=\sum_{i=0}^{N_{j}} \beta_{j, i} e_{i}^{*}
$$

belong to $S^{*}$ and satisfy

$$
\begin{equation*}
\left\{x \in S: f_{j}(x)>\delta_{j}^{2}\right\} \subset\left\{x \in S: g_{j}(x)>\epsilon_{j}\right\} \subset\left\{x \in S: g_{j}(x)>\gamma_{j}\right\} \subset\left\{x \in S: f_{j}(x)>\delta_{j}^{4}\right\} \tag{4.1}
\end{equation*}
$$ for all $j=2, \ldots, k$. Assume also that we have defined numbers $\gamma_{j}$ and functions

$$
h_{j}=\theta_{j}\left(g_{j}\right) \prod_{i<j}\left(1-\theta_{i}\left(g_{i}\right)\right)
$$

where $\theta_{j}: \mathbb{R} \rightarrow[0,1]$ are $C^{\infty}$ functions satisfying

$$
\begin{aligned}
\theta_{j}(t)=0 & \text { if and only } t \leq \gamma_{j} \\
\theta_{j}(t)=1 & \text { if and only } t \geq \epsilon_{j}
\end{aligned}
$$

The interior of the support of $h_{j}$ is the set

$$
U_{j}=\left\{x \in S^{+}: g_{1}(x)<\epsilon_{1}, \ldots, g_{j-1}(x)<\epsilon_{j-1} \text { and } g_{j}(x)>\gamma_{j}\right\}
$$

Then we may again use the density of $\operatorname{span}\left\{e_{n}^{*}: n \in \mathbb{N}\right\}$ in $E^{*}$, in order to find a positive integer $N_{k+1} \geq N_{k}$, and constants $\gamma_{k+1}, \epsilon_{k+1} \in(0,1)$, and $\beta_{k+1,0}, \ldots, \beta_{k+1, N_{k+1}} \in \mathbb{R}$ so that, for

$$
g_{k+1}:=\sum_{j=0}^{N_{k+1}} \beta_{k+1, j} e_{j}^{*}
$$

we have that $g_{k+1} \in S^{*}$ and

$$
\begin{aligned}
& \left\{x \in S: f_{k+1}(x)>\delta_{k+1}^{2}\right\} \subset\left\{x \in S: g_{k+1}(x)>\epsilon_{k+1}\right\} \subset \\
& \left\{x \in S: g_{k+1}(x)>\gamma_{k+1}\right\} \subset\left\{x \in S: f_{k+1}(x)>\delta_{k+1}^{3}\right\} .
\end{aligned}
$$

We now set

$$
\begin{equation*}
U_{k+1}:=\left\{x \in S^{+}: g_{1}(x)<\epsilon_{1}, \ldots, g_{k}(x)<\epsilon_{k} \text { and } g_{k+1}(x)>\gamma_{k+1}\right\} \tag{4.2}
\end{equation*}
$$

and define

$$
\begin{aligned}
& h_{k+1}: S^{+} \longrightarrow \mathbb{R} \\
& h_{k+1}=\theta_{k+1}\left(g_{k+1}\right) \prod_{j<k+1}\left(1-\theta_{j}\left(g_{j}\right)\right) .
\end{aligned}
$$

where $\theta_{k+1}: \mathbb{R} \rightarrow[0,1]$ is a $C^{\infty}$ function such that

$$
\begin{array}{ll}
\theta_{k+1}(t)=0 & \text { if and only if } t \leq \gamma_{k+1} \\
\theta_{k+1}(t)=1 & \text { if and only if } t \geq \epsilon_{k+1} .
\end{array}
$$

Clearly the interior of the support of $h_{k+1}$ is the set $U_{k+1}$.
Thus a sequence $\left\{h_{n}\right\}_{n \in \mathbb{N}}$ of $C^{1}$ smooth functions with the above properties is well defined by induction. As in the proof of Theorem 1.6 it is not difficult to check that the family $\left\{U_{k}\right\}_{k \in \mathbb{N}}$ is a locally finite open covering of $S^{+}$refining $\left\{P_{k}\right\}_{k \in \mathbb{N}}$. Therefore the functions

$$
\psi_{n}:=\frac{h_{n}}{\sum_{k=1}^{\infty} h_{k}}, n \in \mathbb{N},
$$

define a $C^{1}$ partition of unity in $S^{+}$subordinate to $\left\{P_{k}\right\}_{k \in \mathbb{N}}$.
We will also need the following fact.
Claim 4.1. For every $k \in \mathbb{N}$ and every $y \in S^{+}$, we have

$$
h_{k}(y)=0 \Longrightarrow D h_{k}(y)=0 .
$$

Proof. Proceed as in the proof of Claim 3.2,
Again we have that for every $x \in S^{+}$there exist a number $n=n_{x}$ and an open neighborhood $V_{x}$ of $x$ in $S^{+}$such that $\psi_{k}(y)=0$ for every $k>n$ and every $y \in V_{x}$. Let us also call $m=m_{y}$ the largest $j$ such that $h_{j}(y) \neq 0$; this is also the largest $j$ for which $\psi_{j}(y) \neq 0$. Thus, for every $y \in V_{x}$, we have

$$
m_{y}=\max \left\{j: y \in \psi_{j}^{-1}((0,1])\right\} \leq n_{x} .
$$

The derivatives of the functions $\psi_{n}$ can be calculated as in the proof of Theorem 1.6. We have

$$
D \psi_{k}(y)(v)=\sum_{j=1}^{n} \lambda_{k, j}(y) g_{j}(v)=\sum_{j=1}^{m_{y}} \lambda_{k, j}(y) g_{j}(v)
$$

for each $v \in T_{y} S^{+}$, where the functions $\lambda_{k, j}: V_{x} \rightarrow \mathbb{R}$ are of class $C^{1}$.
The following lemma summarizes the properties of the partition of unity $\left\{\psi_{n}\right\}_{n \in \mathbb{N}}$ which will be most useful to us.

Lemma 4.2. Given two continuous functions $f: S^{+} \rightarrow F$ and $\varepsilon: S^{+} \rightarrow(0, \infty)$, there exists a collection of norm-one linear functionals $\left\{g_{k}\right\}_{k \in \mathbb{N}} \subset Y^{*}$ of the form

$$
g_{k}=\sum_{j=0}^{N_{k}} \beta_{k, j} e_{j}^{*},
$$

where $N_{1} \leq N_{2} \leq N_{3} \leq \ldots$, an open covering $\left\{P_{n}\right\}_{n \in \mathbb{N}}$ of $S^{+}$, and a $C^{1}$ partition of unity $\left\{\psi_{n}\right\}_{n \in \mathbb{N}}$ in $S^{+}$such that:
(1) $\left\{\psi_{n}\right\}_{n \in \mathbb{N}}$ is subordinate to $\left\{P_{n}\right\}_{n \in \mathbb{N}}$.
(2) For every $x \in S^{+}$there exist a neighborhood $V_{x}$ of $x$ in $S^{+}$and a number $n=n_{x} \in \mathbb{N}$ such that $\psi_{m}=0$ on $V_{x}$ for all $m>n$, and the derivatives of the functions $\psi_{1}, \ldots, \psi_{n}$ on $V_{x}$ are of the form

$$
D \psi_{k}(y)(v)=\sum_{j=1}^{n} \lambda_{k, j}(y) g_{j}(v)=\sum_{j=1}^{m_{y}} \lambda_{k, j}(y) g_{j}(v)
$$

for $v \in T_{y} S^{+}$, the tangent hyperplane to $S^{+}$at $y \in S^{+} \cap V_{x}$, where $m_{y} \leq n$ is the largest number such that $\psi_{m_{y}}(y) \neq 0$. More precisely, if $L_{y}$ denotes the derivative of the function $u \mapsto \sqrt{1-\|u\|^{2}}$ evaluated at the point $u_{y}$ such that $y=\left(u_{y}, \sqrt{1-\left\|u_{y}\right\|^{2}}\right)$, we have

$$
D \psi_{k}(y)(v)=D \psi_{k}(y)\left(u, L_{y}(u)\right)=L_{y}(u) \mu_{k, 0}(y)+\sum_{j=1}^{N_{n}} \mu_{k, j}(y) e_{j}^{*}(u)=L_{y}(u) \mu_{k, 0}(y)+\sum_{j=1}^{N_{m_{y}}} \mu_{k, j}(y) e_{j}^{*}(u)
$$

for every $k=1, \ldots, n$, and for every $v=\left(u, L_{y}(u)\right) \in T_{y} S^{+}$, where the functions $\mu_{k, j}: V_{x} \rightarrow \mathbb{R}$ are of class $C^{1}, j=1, \ldots, n$.
(3) For every $n \in \mathbb{N}$ there exist a point $y_{n}:=x_{n} \in P_{n}$ such that

$$
\frac{15}{16} \varepsilon\left(y_{n}\right) \leq \varepsilon(y) \leq \frac{17}{16} \varepsilon\left(y_{n}\right) \quad \text { for every } y \in P_{n}, \quad \text { and } \quad \sup _{x, y \in P_{n}}\|f(x)-f(y)\| \leq \frac{\varepsilon\left(y_{n}\right)}{16}
$$

Note also that the integer $n_{x}$ can be chosen as the first such $j$ so that $x \in\left\{y \in S^{+}: g_{j}(y)>\epsilon_{j}\right\}$. Then $V_{x}$ can be chosen as $\left\{y \in S^{+}: g_{j}(y)>\epsilon_{j}\right\}$ and, hence, we have $R_{j} \subset V_{x} \subset P_{j}$ for such $V_{x}$.
Now we are ready to start the construction of our approximating function $\varphi: S^{+} \rightarrow F$, which will be of the form

$$
\varphi(x)=\sum_{n=1}^{\infty}\left(f\left(y_{n}\right)+T_{n}(x)\right) \psi_{n}(x)
$$

where the $y_{n}$ are the points given by condition (3) of the preceding lemma, and the operators $T_{n}$ : $Y \rightarrow F$ will be defined below. We have to distinguish two cases.
Case 1: Assume that $F$ is infinite-dimensional.
In order to define the operators $T_{n}$, we work with the infinite subset $\mathbb{P}$ of $\mathbb{N}$ given by assumption (3) of Theorem 1.7, and we take a countable pairwise disjoint family of infinite subsets of $\mathbb{P}$ which goes to infinity. More precisely, we write

$$
\bigcup_{n=1}^{\infty} I_{n} \subseteq \mathbb{P}
$$

in such a way that:
(1) $I_{n}:=\left\{n_{i}: i \in \mathbb{N}\right\}$ is infinite for each $n \in \mathbb{N}$;
(2) $I_{n} \cap I_{m}=\emptyset$ for all $n \neq m$; and
(3) $\left\{1, \ldots, N_{n}\right\} \cap I_{n}=\emptyset$ for all $n \in \mathbb{N}$.

Here $\left\{N_{n}\right\}_{n \in \mathbb{N}}$ is the non-decreasing sequence of positive integers that appears in the construction of the functionals $g_{n}$ of Lemma 4.2.
Now, by using assumption (3) of the statement, we can find, for each number $n \in \mathbb{N}$, a linear continuous surjection $S_{n}: E \rightarrow F$ of the form

$$
S_{n}=A_{n} \circ P_{n},
$$

where $A_{n}$ is a bounded linear operator from $\overline{\operatorname{span}}\left\{e_{n_{k}}: k \in \mathbb{N}\right\}=\overline{\operatorname{span}}\left\{e_{m}: m \in I_{n}\right\}$ onto $F$, and $P_{n}: E \rightarrow \overline{\operatorname{span}}\left\{e_{n_{k}}: k \in \mathbb{N}\right\}$ is the natural projection associated to the unconditional basis $\left\{e_{j}\right\}_{j \in \mathbb{N}}$.
Now we finally define $T_{n}: Y \rightarrow F$ by

$$
T_{n}(u, t)=\frac{\varepsilon\left(y_{n}\right)}{4\left\|S_{n}\right\|} S_{n}(u)
$$

and $\varphi: S^{+} \rightarrow F$ by

$$
\varphi(x)=\sum_{n=1}^{\infty}\left(f\left(y_{n}\right)+T_{n}(x)\right) \psi_{n}(x) .
$$

It is clear that $\varphi$ is well defined and of class $C^{1}$.
Claim 4.3. We have that $\|\varphi(x)-f(x)\| \leq \varepsilon(x)$ for every $x \in S^{+}$.
Proof. This is shown exactly as in Claim 3.5.
Let us now calculate the derivative of our function $\varphi$. For every $v=\left(u, L_{y}(u)\right) \in T_{y} S^{+}$we have that

$$
\begin{equation*}
T_{n}(v)=T_{n}\left(u, L_{y}(u)\right)=\frac{\varepsilon\left(y_{n}\right)}{4\left\|S_{n}\right\|} S_{n}(u), \tag{4.3}
\end{equation*}
$$

and we have that $D\left(T_{\left.n\right|_{S^{+}}}\right)(y)$ is the restriction of $D T_{n}(y)=T_{n}$ to $T_{y} S^{+}$, that is to say, if $v=$ $\left(u, L_{y}(u)\right) \in T_{y} S^{+}$then

$$
\begin{equation*}
D T_{n}(y)\left(u, L_{y}(u)\right)=\frac{\varepsilon\left(y_{n}\right)}{4\left\|S_{n}\right\|} S_{n}(u) . \tag{4.4}
\end{equation*}
$$

We can now compute the derivative of $\varphi$ on $S^{+}$. Recall that, by condition (2) of Lemma 4.2, for every $x \in S^{+}$there is a neighborhood $V_{x}$ of $x$ in $S^{+}$and a number $n=n_{x} \in \mathbb{N}$ such that

$$
\varphi(y)=\sum_{j=1}^{n}\left(f\left(y_{j}\right)+T_{j}(y)\right) \psi_{j}(y)
$$

for every $y \in V_{x}$. Fix $y \in V_{x}$ and recall that for $m=m_{y}$ (the largest number $j$ for which $\psi_{j}(y) \neq 0$ ), we have

$$
\varphi(y)=\sum_{j=1}^{m}\left(f\left(y_{j}\right)+T_{j}(y)\right) \psi_{j}(y)
$$

By using (4.4) and the expression for $D \psi_{j}(y)$ given in Lemma 4.2, we see that

$$
\begin{align*}
D \varphi(y)\left(u, L_{y}(u)\right) & =\left(\sum_{j=1}^{n} \psi_{j}(y) \frac{\varepsilon\left(y_{j}\right)}{4\left\|S_{j}\right\|} S_{j}(u)\right)+L_{y}(u) \alpha_{N_{n}, 0}(y)+\sum_{j=1}^{N_{n}} \alpha_{N_{n}, j}(y) e_{j}^{*}(u)= \\
& =\left(\sum_{j=1}^{m} \psi_{j}(y) \frac{\varepsilon\left(y_{j}\right)}{4\left\|S_{j}\right\|} S_{j}(u)\right)+L_{y}(u) \alpha_{N_{m}, 0}(y)+\sum_{j=1}^{N_{m}} \alpha_{N_{m}, j}(y) e_{j}^{*}(u) \tag{4.5}
\end{align*}
$$

for every $\left(u, L_{y}(u)\right) \in T_{y} S^{+}, y \in V_{x}$, where the functions $\alpha_{N_{n}, j}: V_{x} \subset S^{+} \rightarrow F$ are of class $C^{1}$ (because we have $\alpha_{N_{n}, j}(y)=\sum_{i=1}^{n}\left(f\left(y_{i}\right)+T_{i}(y)\right) \mu_{i, j}(y)$ and $\alpha_{N_{n}, 0}(y)=\sum_{i=1}^{n}\left(f\left(y_{i}\right)+T_{i}(y)\right) \mu_{i, 0}(y)$, where $\mu_{i, j}$ are as in Lemma 4.2).
Let us now prove that the critical set of $\varphi$ is relatively small.
Lemma 4.4. The set $C_{\varphi}:=\left\{x \in S^{+}: D \varphi(x)\right.$ is not surjective $\}$ is of the form

$$
C_{\varphi}=\left\{\left(w, \sqrt{1-\|w\|^{2}}\right): w \in A\right\}
$$

where $A \subset E$ is a relatively closed subset of the open unit ball of $E$ that is locally contained in a complemented subspace of infinite codimension in $E$.
Proof. Observe that if $n=n_{x}=1$ then $\varphi(y)=f\left(y_{1}\right)+T_{1}(y)$ for every $y \in V_{x}$, and because $T_{1}$ is surjective $\varphi$ does not have any critical point in $V_{x}$. Now let us assume that $n=n_{x} \geq 2$. Let $y=\left(w, \sqrt{1-\|w\|^{2}}\right)$ be point of $V_{x}$. Let $m=m_{y}$ be the largest number such that $\psi_{m}(y) \neq 0$. Recall that $m \leq n$. We only need to show that if

$$
w \notin \overline{\operatorname{span}}\left\{e_{j}: j \in \mathbb{P} \text { or } j=1, \ldots, N_{n}\right\}
$$

then for every $v \in F$ there exists $u \in E$ such that

$$
D \varphi\left(w, \sqrt{1-\|w\|^{2}}\right)\left(u, L_{y}(u)\right)=v
$$

since this will mean that the set

$$
A:=\left\{w \in E:\left(w, \sqrt{1-\|w\|^{2}}\right) \in C_{\varphi}\right\}
$$

will be locally contained in subspaces of the form $\overline{\operatorname{span}}\left\{e_{i}: i \in \mathbb{P}\right.$ or $\left.i=1, \ldots, N_{n}\right\}$, which are complemented, and of infinite codimension, in $E$.
We will need to use the following.
Fact 4.5. For every $w=\sum_{j=1}^{\infty} w_{j} e_{j} \in E \backslash\{0\}$ and every $j_{0} \in \mathbb{N}$ we have that

$$
w_{j_{0}} \neq 0 \Longrightarrow\left\langle J(w), e_{j_{0}}\right\rangle \neq 0
$$

where $J(w)$ denotes $D\|\cdot\|(w)$, and $\langle J(w), u\rangle:=J(w)(u)$.
Proof. If $w_{j_{0}} \neq 0$ then, by assumption (2) of the statement of Theorem 1.7, we have that

$$
\left\|\sum_{j=1, j \neq j_{0}}^{\infty} w_{j} e_{j}\right\| \leq\left\|\sum_{j=1}^{\infty} w_{j} e_{j}\right\| .
$$

This means that the convex function $\theta: \mathbb{R} \rightarrow \mathbb{R}$ defined by

$$
\theta(t)=\left\|w+t e_{j_{0}}\right\|
$$

has a minimum at $t=-w_{j_{0}}$. On the other hand, if we had $\left\langle J(w), e_{j_{0}}\right\rangle=0$, then the same function $\theta$ would have another minimum at the point $t=0$. But since $\|\cdot\|$ is strictly convex the function $\theta$ can only attain its minimum at a unique point. Therefore we must have $\left\langle J(w), e_{j_{0}}\right\rangle \neq 0$.
So let us pick a point $w \in E \backslash \overline{\operatorname{span}}\left\{e_{j}: j \in \mathbb{P}\right.$ or $\left.j=1, \ldots, N_{n}\right\}$ and a vector $v \in F$, and let us construct a vector $u \in E$ such that $D \varphi\left(w, \sqrt{1-\|w\|^{2}}\right)\left(u, L_{y}(u)\right)=v$, where $y=\left(w, \sqrt{1-\|w\|^{2}}\right)$. By assumption, there exists $j_{0} \in \mathbb{N}, j_{0} \in \mathbb{N} \backslash \mathbb{P}$, such that $j_{0}>N_{n} \geq N_{m}$ and $w_{j_{0}} \neq 0$. According to the fact just shown, we have $\left\langle J(w), e_{j_{0}}\right\rangle \neq 0$. Now, since $S_{m}$ is surjective and $\psi_{m}(y) \neq 0$, we may find a sequence $\left(u_{m_{i}}\right)_{i \in \mathbb{N}}$ (indexed by the subsequence $\left(m_{i}\right)_{i \in \mathbb{N}}$ defined by $\left.I_{m}\right)$ such that

$$
\psi_{m}(y) \frac{\varepsilon\left(y_{m}\right)}{4\left\|S_{m}\right\|} S_{m}\left(\sum_{i=1}^{\infty} u_{m_{i}} e_{m_{i}}\right)=v .
$$

Note that $j_{0} \notin I_{m}=\left\{m_{i}: i \in \mathbb{N}\right\}$, because $j_{0} \in \mathbb{N} \backslash \mathbb{P}$ and $I_{m} \subset \mathbb{P}$. Then we can set

$$
u_{j_{0}}:=-\frac{\left\langle J(w), \sum_{i=1}^{\infty} u_{m_{i}} e_{m_{i}}\right\rangle}{\left\langle J(w), e_{j_{0}}\right\rangle}
$$

so that we have

$$
\left\langle J(w), u_{j_{0}} e_{j_{0}}+\sum_{i=1}^{\infty} u_{m_{i}} e_{m_{i}}\right\rangle=0
$$

which bearing in mind that

$$
L_{y}=-\frac{\|w\|}{\sqrt{1-\|w\|^{2}}}\langle J(w), \cdot\rangle
$$

also implies that

$$
L_{y}\left(u_{j_{0}} e_{j_{0}}+\sum_{i=1}^{\infty} u_{m_{i}} e_{m_{i}}\right)=0
$$

So if we set $u_{j}=0$ for all $j \notin I_{m} \cup\left\{j_{0}\right\}$ and we define

$$
u:=\sum_{j=1}^{\infty} u_{j} e_{j}
$$

then we have that

$$
\psi_{m}(y) \frac{\varepsilon\left(y_{m}\right)}{4\left\|S_{m}\right\|} S_{m}(u)=v, \quad L_{y}(u)=0, \sum_{j=1}^{N_{m}} \alpha_{N_{m}, j}(y) e_{j}^{*}(u)=0, \quad \text { and also } S_{j}(u)=0 \text { for } j<m
$$

because $j_{0}>N_{n} \geq N_{m}, I_{m} \cap\left\{1,2, \ldots, N_{m}\right\}=\emptyset$, and the sets $I_{j}$ are pairwise disjoint. In view of (4.5)) these equalities imply that $D \varphi(y)(u)=v$.
Now, according to Lemma 4.4 and Theorem 1.4, we can extract the set $C_{\varphi}$, since it is $C^{1}$ diffeomorphic (via the projection of the graph $S^{+}$of the function $w \mapsto \sqrt{1-\|w\|^{2}}$ onto the open unit ball of $E$ ) to a subset which can be extracted. Therefore we can finish the proof of Theorem 1.7 exactly as we did with Theorem 1.6,
Case 2: Assume that $F=\mathbb{R}^{m}$. The proof is almost identical, but with the following important difference: now the set $\mathbb{P}$ is by definition the set of even positive integers, and the sets $I_{n}$ are finite subsets of $\mathbb{P}$ such that:
(1) $\sharp I_{n}=m$ for each $n \in \mathbb{N}$;
(2) $I_{n} \cap I_{j}=\emptyset$ for all $n \neq j$; and
(3) $\left\{1, \ldots, N_{n}\right\} \cap I_{n}=\emptyset$ for all $n \in \mathbb{N}$.

Here $\left\{N_{n}\right\}_{n \in \mathbb{N}}$ is the non-decreasing sequence of positive integers that appears in the construction of the functionals $g_{n}$ of Lemma 4.2.
Of course in this case we can always find linear surjections $A_{n}: \operatorname{span}\left\{e_{i}: i \in I_{n}\right\} \rightarrow \mathbb{R}^{m}$.

## 5. Technical versions of Theorems 1.6 and 1.7, examples, and remarks

In this section we will give some examples, make some remarks and establish more technical variants of our results which follow by the same method of proof. We will also prove Proposition 1.8,
The proof of Theorem 1.6 can be easily adjusted to obtain more general results with more complicated statements. Namely, the following two results are true.
Theorem 5.1. Let $E$ and $F$ be Banach spaces. Assume that:
(1) $E$ is infinite-dimensional, with a separable dual $E^{*}$.
(2) There exist three sequences $\left\{E_{n, 1}\right\}_{n \geq 1},\left\{E_{n, 2}\right\}_{n \geq 1},\left\{E_{n, 3}\right\}_{n \geq 2}$ of subspaces of $E$ such that $E=E_{1,1} \oplus E_{1,2}, \quad E_{1,2}=\left(E_{2,1} \oplus E_{2,2}\right) \oplus \ldots \oplus\left(E_{n, 1} \oplus E_{n, 2} \oplus E_{n, 3}\right), \quad E_{n, 3}=E_{n+1,1} \oplus E_{n+1,2} \oplus E_{n+1,3}$,
with either $E_{n, 3}$ infinite-dimensional and reflexive and $\operatorname{dim} E_{n, 2} \geq 1$, or else $E_{n, 2}$ infinitedimensional and reflexive for all $n \geq 2$. Suppose also that there exists a bounded linear operator from $E_{n, 1}$ onto $F$ for every $n \in \mathbb{N}$.
Then, for every continuous mapping $f: E \rightarrow F$ and every continuous function $\varepsilon: E \rightarrow(0, \infty)$ there exists a $C^{1}$ mapping $g: E \rightarrow F$ such that $\|f(x)-g(x)\| \leq \varepsilon(x)$ and $D g(x): E \rightarrow F$ is a surjective linear operator for every $x \in E$.

Observe that if $E_{n, 3}$ is infinite-dimensional and reflexive, the spaces $E_{n, 2}$ can be taken to be of dimension 1 for every $n \in \mathbb{N}$. The proof is almost the same as that of Theorem 1.6 . Here, up to finite-dimensional perturbations of the subspaces $E_{k, j}$, we can arrange that $E_{1,2} \subset K e r g_{1}^{1}$ and that $E_{n, 1} \oplus E_{n, 3} \subseteq \bigcap_{j=1}^{n} \operatorname{Kerg}{ }_{j}^{1}$, and we may set $A_{x}=\left\{y \in S^{+}: E_{n, 3} \subset \operatorname{Ker} L_{y}\right\}$.
Theorem 5.2. Let $E, X$, and $F$ be Banach spaces. Assume either that $E$ is infinite-dimensional, separable, and reflexive, and $F$ is finite-dimensional, or that:
(1) $E$ is infinite-dimensional, with a separable dual $E^{*}$.
(2) There exists a decomposition of $E$,

$$
E=G \oplus E_{1} \oplus X
$$

such that $G$ is infinite-dimensional and reflexive, and $E_{1}$ is isomorphic to $E$.
(3) There exists a bounded linear operator from $G \oplus X$ onto $F$ (equivalently, $F$ is a quotient of $G \oplus X)$.
Then, for every continuous mapping $f: E \rightarrow F$ and every continuous function $\varepsilon: E \rightarrow(0, \infty)$ there exists a $C^{1}$ mapping $g: E \rightarrow F$ such that $\|f(x)-g(x)\| \leq \varepsilon(x)$ and $D g(x): E \rightarrow F$ is a surjective linear operator for every $x \in E$.

If, additionally, $X$ is isomorphic to $E$, then $F$ can be taken as a quotient of $E$.
Observe that each of these two results imply Theorem 1.6 , with Theorem 5.1 being the most general one.
For instance, Theorem 5.2 can be applied to the James space $J$ and to its dual $J^{*}$. Indeed, both spaces have separable dual. It is known that $J$ has many reflexive infinite-dimensional complemented subspaces $G$ [22]. Since $J$ is prime [21], for each such $G$, we can write $J=G \oplus J$ (for instance we have $\left.J=l_{2} \oplus J\right)$. Now, recalling the fact that $J$ has a separable dual, apply Theorem 5.2 to $E=E_{1}=J$ and $X=\{0\}$ to see that every continuous function $f: J=G \oplus J \rightarrow F$, where $F$ is a quotient of $G$, can be uniformly approximated by $C^{1}$ smooth mappings without critical points. Similar arguments work for the dual of the James space $J^{*}$.

It also follows that the conclusion of Theorem 1.6 is true for composite spaces of the form $c_{0} \oplus \ell_{p}$ or $c_{0} \oplus L^{p}, 1<p<\infty$, with $k$ being the order of smoothness of $\ell_{p}$ or $L^{p}$. More generally, if $E$ is any finite direct sum of the classical Banach spaces $c_{0}, \ell_{p}$ or $L^{p}, 1<p<\infty$, and there is a bounded linear operator from $E$ onto $F$ then the conclusion of Theorem 1.2 is true with $k$ being the minimum of the orders of smoothness of the spaces appearing in this decomposition of $E$.

Remark 5.3. Notice that that in the case that $E$ is a separable Hilbert space, we have that the function $w \mapsto\|w\|^{2}$ is of class $C^{\infty}$, hence all the mappings appearing in the proof of Theorem 1.6 are of class $C^{\infty}$, and we directly obtain an approximating function $g$ of class $C^{\infty}$ with no critical points.

In fact, in the Hilbertian case we do not need to use a partition of unity in the upper sphere $S^{+}$. We can directly construct a partition of unity $\left\{\psi_{n}\right\}_{n \in \mathbb{N}}$ in $E$ subordinated to an open covering by open balls with linearly independent centers $\left\{y_{j}\right\}$, as in [6]. Then, choosing an orthonormal basis $\left\{e_{j}\right\}$ for which $\operatorname{span}\left\{y_{1}, \ldots, y_{n}\right\}=\operatorname{span}\left\{e_{1}, \ldots, e_{n}\right\}$ for every $n \in \mathbb{N}$, we define operators $T_{n}: E \rightarrow F$ as in the proof of Theorem 1.7, where $\mathbb{P}$ can be any infinite subset of $\mathbb{N}$ such that $\mathbb{N} \backslash \mathbb{P}$ is also infinite. Then one can easily check that the function

$$
\varphi(y)=\sum_{n=1}^{\infty}\left(f\left(y_{n}\right)+T_{n}\left(y-y_{n}\right)\right) \psi_{n}(y)
$$

approximates $f$ and the set $C_{\varphi}$ of its critical points is locally contained in a subspace of infinite codimension in $E$, specifically in subspaces of the form $\overline{\operatorname{span}}\left\{e_{j}: j \in \mathbb{P}\right.$ or $\left.j=1, \ldots, n\right\}$. Then one can extract $C_{\varphi}$ by means of a $C^{\infty}$ diffeomorphism $h: E \rightarrow E \backslash C_{\varphi}$ which is sufficiently close to the identity, and conclude that the function $g:=\varphi \circ h$ approximates $f$ and has no critical points.

The same proof as that of Theorem 1.7, with some adjustments, allows us to obtain a more general (and also more technical) result as follows.

Theorem 5.4. Let $E$ be an infinite-dimensional Banach space, and $F$ be a Banach space such that:
(1) E has an equivalent norm $\|\cdot\|$ which is $C^{1}$ and locally uniformly convex.
(2) E has a (normalized) Schauder basis $\left\{e_{n}\right\}_{n \in \mathbb{N}}$ which is shrinking.
(3) There exists an infinite subset $\mathbb{I}$ of $\mathbb{N}$ such that the subspace $\overline{\operatorname{span}}\left\{e_{j}: j \in \mathbb{N} \backslash \mathbb{I}\right\}$ is complemented in $E$, and for every $x=\sum_{j=}^{\infty} x_{j} e_{j}$ and every $j_{0} \in \mathbb{I}$ we have that

$$
\left\|\sum_{j \in \mathbb{N}, j \neq j_{0}} x_{j} e_{j}\right\| \leq\left\|\sum_{j \in \mathbb{N}} x_{j} e_{j}\right\| .
$$

(4) In the case that $F$ is infinite-dimensional, there exists an infinite subset $\mathbb{P}$ of $\mathbb{N}$ such that $\mathbb{I} \backslash \mathbb{P}$ is infinite and for every infinite subset $J$ of $\mathbb{P}$ the subspace $E^{\prime}=\overline{\operatorname{span}}\left\{e_{j}: j \in J \cup(\mathbb{N} \backslash \mathbb{I})\right\}$ is complemented in $E$, and there exists a linear bounded operator from $E^{\prime}$ onto $F$.
Then, for every continuous mapping $f: E \rightarrow F$ and for every continuous function $\varepsilon: E \rightarrow(0, \infty)$ there exists a $C^{1}$ mapping $g: E \rightarrow F$ such that $\|f(x)-g(x)\| \leq \varepsilon(x)$ and $D g(x): E \rightarrow F$ is a surjective linear operator for every $x \in E$.

Proof. The only important difference with the proof of Theorem 1.7 is that now we have to use an analogue of Fact 4.5 which is true if we just pick $j_{0} \in \mathbb{I}$. Therefore, if we take $w=\sum_{j=1}^{\infty} w_{j} e_{j} \notin$ $\overline{\operatorname{span}}\left\{e_{j}: j \in \mathbb{P} \cup(\mathbb{N} \backslash \mathbb{I})\right.$ or $\left.j=1, \ldots, N_{n}\right\}$, since $\mathbb{I} \backslash \mathbb{P}$ is infinite, there will exist $j_{0} \in \mathbb{I}, j_{0}>N_{n}$ such that $w_{j_{0}} \neq 0$ and thus $\left\langle J(w), e_{j_{0}}\right\rangle \neq 0$. The operators $S_{n}$ have supports in complemented subspaces of the form $\overline{\operatorname{span}}\left\{e_{j}: j \in I_{n} \cup(\mathbb{N} \backslash \mathbb{I})\right\}$, where the sets $I_{n} \subset \mathbb{P}$ are defined as in the proof of Theorem 1.7.

Remark 5.5. It is clear that the spaces $\ell_{p}$ and $L^{p}, 1<p<\infty$ satisfy the assumptions of Theorem 1.6. It may not be so obvious why the space $c_{0}$ satisfy the assumptions of Theorem 1.7, let us clarify this point. If we repeat the proof of [26, Theorem V.1.5] in the particular case that $\Gamma=\mathbb{N}$, since all the operations that are made in this proof are coordinate-wise monotone, we see that the $C^{1}$ and LUR renorming $\|\cdot\|$ that we obtain for $c_{0}$ has the property that

$$
\left\|\sum_{j \in \mathbb{N}, j \neq j_{0}} x_{j} e_{j}\right\| \leq\left\|\sum_{j \in \mathbb{N}} x_{j} e_{j}\right\|
$$

for every $j_{0} \in \mathbb{N}$ and every $x=\left(x_{1}, x_{2}, x_{3}, \ldots\right) \in c_{0}$, where $\left\{e_{n}\right\}$ is the canonical basis of $c_{0}$. This shows that this norm $\|\cdot\|$ satisfies assumptions (1) and (2) of Theorem 1.7. On the other hand, for every infinite subset $J$ of $\mathbb{N}$ we have that $\overline{\operatorname{span}}\left\{e_{j}: j \in J\right\}$ is isomorphic to $c_{0}$, so it is clear that assumption (3) is satisfied as well, provided that there exists a continuous linear operator from $c_{0}$ onto $F$. Therefore $E=c_{0}$ satisfies the conclusion of Theorem 1.2,

The latter fact can be generalized to Banach spaces with a shrinking basis which contain copies of $c_{0}$.
Theorem 5.6. Let $E$ be a Banach space that contains the space $c_{0}$ and admits a shrinking Schauder basis. Let $F$ be a quotient of $E$.
Then, for every continuous mapping $f: E \rightarrow F$ and for every continuous function $\varepsilon: E \rightarrow(0, \infty)$ there exists a $C^{1}$ mapping $g: E \rightarrow F$ such that $\|f(x)-g(x)\| \leq \varepsilon(x)$ and $D g(x): E \rightarrow F$ is a surjective linear operator for every $x \in E$.

Proof. We will show that $E$ satisfies the assumptions of Theorem 5.4.
First, by Sobczyk's Theorem [64], $c_{0}$ is complemented in $E$, that is, $E$ is isomorphic to $G \oplus c_{0}$, for a certain Banach space $G$. Since $c_{0}$ is isomorphic to $c_{0} \oplus c_{0}, G$ may be taken as $E$. So, we can and will assume that

$$
E=c_{0} \oplus E .
$$

Let $\left\{e_{j}\right\}_{j \in N}$ be the canonical Schauder basis in $c_{0}$. Equip $c_{0}$ with the $C^{1}$ and LUR norm $\|\cdot\|$ which was described in Remark [5.5. That is, for every $j_{0} \in \mathbb{N}$, we have

$$
\left\|\sum_{j=1, j \neq j_{0}}^{\infty} \alpha_{j} e_{j}\right\| \leq\left\|\sum_{j=1}^{\infty} \alpha_{j} e_{j}\right\|,
$$

for every $x=\sum_{j=1}^{\infty} \alpha_{j} e_{j} \in c_{0}$.
Similarly, let $\left\{d_{n}\right\}_{n \in \mathbb{N}}$ be a shrinking Schauder basis in $E$. Equip $E$ with a $C^{1}$ and LUR norm $|\cdot|$. Define a new norm in $c_{0} \oplus E$, by letting

$$
\||x+y|\|=\sqrt{\|x\|^{2}+|y|^{2}},
$$

for every $x+y \in c_{0} \oplus E$. This norm is $C^{1}$ and LUR as well. Define $\left\{f_{k}\right\}_{k \in \mathbb{N}} \subset c_{0} \oplus E$, where $f_{2 j-1}=e_{j}+0$ and $f_{2 n}=0+d_{n}$ for every $j, n \in \mathbb{N}$. It is also easy to check that $\left\{f_{k}\right\}_{k \in \mathbb{N}}$ is a shrinking Schauder basis for $c_{0} \oplus E$.
For $x+y \in c_{0} \oplus E$, let $x=\sum_{j=1}^{\infty} \alpha_{j} e_{j} \in c_{0}$ and $y=\sum_{n=1}^{\infty} \beta_{n} d_{n} \in E$ be their basis expansions. Then, writing $z_{2 j-1}=\alpha_{j}$ and $z_{2 n}=\beta_{n}$, we obtain the expansion of $z=x+y=\sum_{k=1}^{\infty} z_{k} f_{k}=$ $\sum_{j=1}^{\infty} z_{2 j-1} e_{j}+\sum_{n=1}^{\infty} z_{2 n} d_{n} \in c_{0} \oplus E$. For every $j_{0} \in \mathbb{N}$, we have

$$
\begin{aligned}
& \left\|\left\|\sum_{k=1, k \neq 2 j_{0}}^{\infty} z_{k} f_{k}\right\|\right\|=\| \| \sum_{j=1, j \neq j_{0}}^{\infty} z_{2 j-1} e_{j}+\sum_{n=1}^{\infty} z_{2 n} d_{n}\| \| \leq\| \| \sum_{j=1, j \neq j_{0}}^{\infty} \alpha_{j} e_{j}+\sum_{n=1}^{\infty} \beta_{n} d_{n}\| \| \leq \\
& \left\|\left\|\left(\sum_{j=1, j \neq j_{0}}^{\infty} \alpha_{j} e_{j}\right)+y\right\|\right\|=\left(\left\|\sum_{j=1, j \neq j_{0}}^{\infty} \alpha_{j} e_{j}\right\|^{2}+|y|^{2}\right)^{\frac{1}{2}} \leq \sqrt{\|x\|^{2}+|y|^{2}}= \\
& \left\|\left|x+y\| \|=\left\|\mid \sum_{k=1}^{\infty} z_{k} f_{k}\right\| \|,\right.\right.
\end{aligned}
$$

where in the second line we have used the fact that $\left\|\sum_{j=1, j \neq j_{0}}^{\infty} \alpha_{j} e_{j}\right\|^{2} \leq\left\|\sum_{j=1}^{\infty} \alpha_{j} e_{j}\right\|^{2}=\|x\|^{2}$. Now, we are in a position to apply Theorem 5.4. Namely, let $\mathbb{I}=\{2 n: n \in \mathbb{N}\}$ and $\mathbb{P}=\{4 n: n \in \mathbb{N}\}$.

Corollary 5.7. Let $C(K)$ be the Banach space of continuous functions, where $K$ is a metrizable countable compactum and $F$ be a quotient of $C(K)$.
Then, for every continuous mapping $f: C(K) \rightarrow F$ and for every continuous function $\varepsilon: C(K) \rightarrow$ $(0, \infty)$ there exists a $C^{\infty}$ mapping $g: E \rightarrow F$ such that $\|f(x)-g(x)\| \leq \varepsilon(x)$ and $D g(x): E \rightarrow F$ is a surjective linear operator for every $x \in E$.

Proof. By an application of [43, Theorem 1.4], which states that a Banach space has a shrinking basis provided its dual has a Schauder basis, we obtain that $C(K)$ has a shrinking basis (because $\left.C(K)^{*}=l_{1}\right)$. Moreover, using the fact that $c_{0}$ is a subspace of $C(K)$, we infer that $C(K)$ is isomorphic to $c_{0} \oplus G$ for some Banach space $G$, which yields (as in the above proof) that $C(K)$ is isomorphic to $c_{0} \oplus C(K)$. Hence, by Theorem 5.6, the $C^{1}$ version of our assertion holds. The $C^{\infty}$ version requires the fact that $C(K)$ has an equivalent $C^{\infty}$ norm, which is due to Haydon [42, and Proposition 1.8 ,

For more information about the spaces $C(K)$ we refer the reader to [ 60 . The space $C(K)$ is an example of isometric predual of $\ell_{1}$ (meaning a Banach space $E$ with an equivalent norm $\|\cdot\|$ such that the dual $\left(E^{*},\|\cdot\|^{*}\right)$ is isometric to $\left.\ell_{1}\right)$. The class of isomorphic predual spaces for $\ell_{1}$ is larger that the class of isometric predual spaces (the space constructed by Bourgain and Delbaen [17] is such an example), which in turn is smaller than the class of $C(K)$ spaces for metrizable countable compactum $K$, see [13].

Remark 5.8. Since every isometric predual space $E$ of $\ell_{1}$ contains $c_{0}$ (see for instance [70, Corollary $1]$ ) and admits an equivalent real-analytic norm [25, Corollary 3.3], the above corollary is valid for $E$. Even more, the corollary is valid for any infinite-dimensional separable Banach space $E$ which has a shrinking basis and which admits an equivalent polyhedral norm (equivalently, with a countable James boundary). This follows from the facts that, being polyhedral, $E$ must contain $c_{0}$, and that a space with a countable James boundary admits an equivalent real-analytic norm (see [25] or [38, Chapter 5, section 6] for reference).

As we noted in the introduction our main results imply that continuous functions between many Banach spaces can be arbitrarily well approximated by smooth open mappings.

Remark 5.9. Let $(E, F)$ be a pair of Banach spaces with the property that for every continuous mapping $f: E \rightarrow F$ and for every continuous function $\varepsilon: E \rightarrow(0, \infty)$ there exists a $C^{k}$ mapping $g: E \rightarrow F$ with no critical points such that $\|f(x)-g(x)\| \leq \varepsilon(x), x \in E$. Then the pair $(E, F)$ also has the following property: for every continuous mapping $f: E \rightarrow F$ and for every continuous function $\varepsilon: E \rightarrow(0, \infty)$ there exists an open mapping $g: E \rightarrow F$ of class $C^{k}$ such that $\|f(x)-g(x)\| \leq \varepsilon(x)$, $x \in E$.

This follows trivially from [49, Theorem XV.3.5]. Recall that $g: E \rightarrow F$ is said to be open if for every open subset $U$ of $E$ we have that $g(U)$ is open in $F$. Notice that the approximation of arbitrary continuous maps by smooth (or even merely continuous) open maps is impossible for $E=\mathbb{R}^{n}$ : for instance, if $E=\mathbb{R}^{n}, F=\mathbb{R}, f(x)=e^{-\|x\|^{2}}, \varepsilon(x)=1 / 3$, every continuous function $g$ which $\varepsilon$ approximates $f$ must attain a global maximum in $\mathbb{R}^{n}$, hence $g\left(\mathbb{R}^{n}\right)$ is not open in $\mathbb{R}$.

Example 5.10. In view of Theorem 1.1 it is perhaps natural to ask whether in the case $E=F$ one can get $C^{\infty}$ approximations $g: E \rightarrow E$ such that $D g(x): E \rightarrow E$ is a linear isomorphism for every $x \in E$. This is not possible, as the following example shows.
Let $f: \ell_{2} \rightarrow \ell_{2}$ be defined by

$$
f\left(\sum_{n=1}^{\infty} x_{n} e_{n}\right)=\left(\sum_{n=1}^{\infty}\left|x_{n}\right| e_{n}\right),
$$

where $\left\{e_{n}\right\}_{n \in \mathbb{N}}$ denotes the usual basis of $\ell_{2}$ (that is, $e_{1}=(1,0,0, \ldots), e_{2}=(0,1,0, \ldots)$, etc). Assume that there exists $g \in C^{\infty}(E, E)$ such that $D g(x): \ell_{2} \rightarrow \ell_{2}$ is an isomorphism and $\|f(x)-g(x)\| \leq 1 / 3$ for every $x \in \ell_{2}$. Consider the projection $P_{1}: \ell_{2} \rightarrow \mathbb{R}$ given by $P_{1}(x)=x_{1}$, and the function $g_{1}=P_{1} \circ g$. Since $D g(x)$ is an isomorphism for every $x$, we must have $D g_{1}(x)=P_{1} \circ D g(x)=D g(x)\left(e_{1}\right) \neq 0$ for every $x \in E$, and in particular, considering the curve $\gamma_{1}(t)=t e_{1}, t \in \mathbb{R}$, and the function

$$
\theta(t):=g_{1}\left(\gamma_{1}(t)\right), \quad t \in \mathbb{R},
$$

we must have

$$
\begin{equation*}
\theta^{\prime}(t)=D g_{1}\left(\gamma_{1}(t)\right)\left(e_{1}\right) \neq 0 \tag{5.1}
\end{equation*}
$$

for all $t \in \mathbb{R}$. However,

$$
\left|P_{1}\left(g\left(\gamma_{1}(t)\right)\right)-P_{1}\left(f\left(\gamma_{1}(t)\right)\right)\right| \leq\left\|g\left(\gamma_{1}(t)\right)-f\left(\gamma_{1}(t)\right)\right\| \leq 1 / 3,
$$

hence

$$
\theta(1)=P_{1}\left(g\left(\gamma_{1}(1)\right)\right) \geq 1-1 / 3=2 / 3,
$$

and similarly

$$
\theta(-1) \geq 2 / 3>1 / 3 \geq \theta(0)
$$

Thus $\theta$ must attain a minimum at some point $t_{0}$ of the interval $(-1,1)$, which implies that $\theta^{\prime}\left(t_{0}\right)=0$ and contradicts (5.1).

Proof of Theorem 1.5. As said in the introduction, by the results of [30, 46], it is sufficient to show Theorem [1.5 for functions $f: U \rightarrow V$, where $U \subset E$ and $V \subset F$ are open subsets of two separable Hilbert spaces $E, F$, respectively. Observe that we can assume $V=F$. Indeed, if $f: U \rightarrow V \subset F$, $\varepsilon: U \rightarrow(0, \infty)$ are continuous functions then, by taking $\widetilde{\varepsilon}(x)=\frac{1}{2} \min \{\varepsilon(x), \operatorname{dist}(f(x), F \backslash V)\}$, if we are able to $\widetilde{\varepsilon}$-approximate $f: U \rightarrow F$ by a smooth function $g: U \rightarrow F$ with no critical points, then we also have that $\|g(x)-f(x)\|<\operatorname{dist}(f(x), F \backslash V)$, which implies that $g(x) \in V$ for every $x \in U$; that is, we really have $g: U \rightarrow V$. On the other hand, showing the result for $f: U \rightarrow F$ is not more difficult than proving it in the case $U=E$ (though it does encumber the notation). For example, it requires a version of the extractibility fact (a counterpart of Theorem 1.3) where the whole space $E$, its closed subset $X$, and an open cover $\mathcal{G}$ of $E$ must be replaced with an open subset $U$ (of $E$ ), a closed subset of $U$, and an open cover of $U$, respectively. Such a fact can be proved by mimicking the technique of the proof of Theorem (1.3); one just has to make some easy adjustments in the appropriate places. We leave the details to the interested reader.

Throughout the paper the "limiting" function $\varepsilon(x)$ is assumed to be positive. The following remark explains what can be said if we merely require that $\varepsilon(x) \geq 0$.

Remark 5.11. Let $H$ be a separable, infinite-dimensional Hilbert space and $f: H \rightarrow H$ be a continuous mapping. Then, for every continuous function $\varepsilon: H \rightarrow[0, \infty)$, there exists a continuous mapping $g: H \rightarrow H$ such that the restriction $g_{\left.\right|_{H \backslash \varepsilon-1}(0)}$ is $C^{\infty}$ smooth and has no critical points, and $\|f(x)-g(x)\| \leq \varepsilon(x)$ for every $x \in H$ (hence, $f(x)=g(x)$ provided $\varepsilon(x)=0$ ). This a consequence of Theorem 1.5 applied to $U=H \backslash \varepsilon^{-1}(0)$ and $\varepsilon_{\left.\right|_{U}}$.

Let us conclude this paper with the proof of Proposition 1.8,
Proof of Proposition 1.8, Let $f: E \rightarrow F$ and $\varepsilon: E \rightarrow(0, \infty)$ be continuous. By assumption (1) there exists a $C^{1}$ function $\varphi: E \rightarrow F$ without critical points so that

$$
\|f(x)-\varphi(x)\| \leq \varepsilon(x) / 2
$$

It is well known that the set of continuous linear surjections from a Banach space $E$ onto a Banach space $F$ is open; see [49, Theorem XV.3.4] for instance. Therefore, for each $x \in E$ there exists $r_{x}>0$ such that if $S: E \rightarrow F$ is a bounded linear operator then

$$
\begin{equation*}
\|S-D \varphi(x)\|<2 r_{x} \Longrightarrow S \text { is surjective. } \tag{5.2}
\end{equation*}
$$

By continuity of $D \varphi$, for every $x$ we may find a number $s_{x} \in\left(0, r_{x}\right)$ such that if $y \in B\left(x, s_{x}\right)$ then

$$
\|D \varphi(y)-D \varphi(x)\|<r_{x} .
$$

Since $E$ is separable, we can extract a countable subcovering

$$
E=\bigcup_{n=1}^{\infty} B\left(x_{n}, s_{n}\right)
$$

where $s_{n}:=s_{x_{n}}$. Let us also denote $r_{n}:=r_{x_{n}}$, and define $\eta: E \rightarrow(0,1)$ by

$$
\eta(y)=\min \left\{\frac{\varepsilon(y)}{2}, \sum_{n=1}^{\infty} \frac{s_{n}}{2} \psi_{n}(y)\right\},
$$

where $\left\{\psi_{n}\right\}$ is a partition of unity such that the open support of $\psi_{n}$ is contained in $B\left(x_{n}, s_{n}\right)$. Now we may apply assumption (2) to find a $C^{k}$ function $g: E \rightarrow F$ such that

$$
\|\varphi(y)-g(y)\| \leq \eta(y), \text { and }\|D \varphi(y)-D g(y)\| \leq \eta(y)
$$

for all $y \in E$. Then for every $y \in E$ there exists $n=n_{y} \in \mathbb{N}$ such that $y \in B\left(x_{n}, s_{n}\right)$ and $\eta(y) \leq s_{n} / 2$. It follows that $\|D g(y)-D \varphi(y)\| \leq s_{n} / 2<r_{n}$ and $\left\|D \varphi(y)-D \varphi\left(x_{n}\right)\right\|<r_{n}$, hence $\left\|D g(y)-D \varphi\left(x_{n}\right)\right\|<$ $2 r_{n}$, and according to (5.2) this implies that $D g(y)$ is surjective. This shows that $g$ has no critical points. On the other hand, since $\eta \leq \varepsilon / 2$, it is clear that

$$
\|f(y)-g(y)\| \leq\|f(y)-\varphi(y)\|+\|\varphi(y)-g(y)\| \leq \varepsilon(y) / 2+\varepsilon(y) / 2=\varepsilon(y)
$$

so $g$ also approximates $f$ as required.

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[^1]:    ${ }^{1}$ For $C^{k}$ with $k \geq 2$ in place of $C^{1}$, the answer to this question is positive; see 23].

[^2]:    ${ }^{2}$ Smoothness of $s(u)$ at $u=0$ is a consequence of the facts that $\|\cdot\|^{2}$ is trivially differentiable at 0 , and that an everywhere differentiable convex function is always of class $C^{1}$; see for instance [16, Corollary 4.2.12].

